

CONFORMALLY INVARIANT SYSTEMS OF DIFFERENTIAL OPERATORS ASSOCIATED TO TWO-STEP NILPOTENT MAXIMAL PARABOLICS OF NON-HEISENBERG TYPE

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ABSTRACT. In this paper we construct conformally invariant systems of first order and second order differential operators associated to a homogeneous line bundle $\mathcal{L}_s \rightarrow G_0/Q_0$ with Q_0 a maximal two-step nilpotent parabolic subgroup of a Lie group G_0 . This generalizes the results by Barchini, Kable, and Zierau. To do so we use techniques different from these used by them.

1. INTRODUCTION

The main work of this paper concerns systems of differential operators that are equivariant under an action of a Lie algebra. We call such systems *conformally invariant*. To explain the meaning of the equivariance condition, suppose that $\mathcal{V} \rightarrow M$ is a vector bundle over a smooth manifold M and \mathfrak{g} is a Lie algebra of first-order differential operators that act on sections of \mathcal{V} . A linearly independent list D_1, \dots, D_n of linear differential operators on sections of \mathcal{V} is called a *conformally invariant system* if, for each $X \in \mathfrak{g}$, there are smooth functions $C_{ij}^X(m)$ on M so that, for all $1 \leq i \leq n$, and sections f of \mathcal{V} , we have

$$(1.1) \quad ([X, D_i] \bullet f)(m) = \sum_{j=1}^n C_{ji}^X(m) (D_j \bullet f)(m),$$

where $[X, D_j] = XD_j - D_jX$. (See Definition 2.3 for the precise definition.) Here, the dot \bullet denotes the action of differential operators on smooth functions.

An important consequence of the definition (1.1) is that the common kernel of the operators in the conformally invariant system D_1, \dots, D_n is invariant under a Lie algebra action. The representation theoretic question of understanding the common kernel as a \mathfrak{g} -module is an open question (except for a small number of very special examples).

The notion of conformally invariant systems generalizes that of quasi-invariant differential operators introduced by Kostant in [14] and is related to a work of Huang ([7]). It is also compatible with the definition given by Ehrenpreis in [5]. Conformally invariant systems are explicitly or implicitly presented in the work of Davidson-Enright-Stanke ([4]), Kable ([9], [10]), Kobayashi-Ørsted ([13]), Wallach ([19]), among others. Much of the published work is for the case that $M = G/Q$ with $Q = LN$, N abelian. The systematic study of conformally invariant systems started with the work of Barchini-Kable-Zierau in [1] and [2].

Although the theory of conformally invariant systems can be viewed as a geometric-analytic theory, it is closely related to algebraic objects such as generalized Verma modules. It has been

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shown in [2] that a conformally invariant system yields a homomorphism between certain generalized Verma modules. The classification of non-standard homomorphisms between generalized Verma modules is an open problem. In [16], it is determined whether or not the homomorphisms between the generalized Verma modules that arise from certain conformally invariant systems are standard.

The main goal of this paper is to build systems of differential operators that satisfy the condition (1.1), when M is a homogeneous manifold G/Q with Q a maximal two-step nilpotent parabolic subgroup. This is to construct systems D_1, \dots, D_n acting on sections of bundles $\mathcal{V}_s \rightarrow G/Q$ over G/Q in a systematic manner and to determine the bundles \mathcal{V}_s on which the systems are conformally invariant. The method that we use is different from one used by Barchini-Kable-Zierau in [1].

To describe our work more precisely, let G be a complex, simple, connected, simply-connected Lie group with Lie algebra \mathfrak{g} . It is known that \mathfrak{g} has a \mathbb{Z} -grading $\mathfrak{g} = \bigoplus_{j=-r}^r \mathfrak{g}(j)$ so that $\mathfrak{q} = \mathfrak{g}(0) \oplus \bigoplus_{j>0} \mathfrak{g}(j) = \mathfrak{l} \oplus \mathfrak{n}$ is a parabolic subalgebra of \mathfrak{g} . Let $Q = N_G(\mathfrak{q}) = LN$. For a real form \mathfrak{g}_0 of \mathfrak{g} , define G_0 to be an analytic subgroup of G with Lie algebra \mathfrak{g} . Set $Q_0 = N_{G_0}(\mathfrak{q})$. Our manifold is $M = G_0/Q_0$ and we consider a line bundle $\mathcal{L}_s \rightarrow G_0/Q_0$ for each $s \in \mathbb{C}$. By the Bruhat theory, that G_0/Q_0 admits an open dense submanifold $\bar{N}_0 Q_0/Q_0$. We restrict our bundle to this submanifold. The systems that we study act on smooth sections of the restricted bundle.

To build systems of differential operators we observe that L acts by the adjoint representation on $\mathfrak{g}(1)$ with a unique open orbit. This makes $\mathfrak{g}(1)$ a prehomogeneous vector space. Our construction is based on the invariant theory of a prehomogeneous vector space. It is natural to associate L -equivariant polynomial maps called *covariant maps* to the prehomogeneous vector space $(L, \text{Ad}, \mathfrak{g}(1))$. To define our systems of differential operators, we use covariant maps $\tau_k : \mathfrak{g}(1) \rightarrow \mathfrak{g}(-r+k) \otimes \mathfrak{g}(r)$ that are associated to $\mathfrak{g}(1)$. (See Definition 3.1.) Each τ_k can be thought of as giving the symbols of the differential operators that we study.

Let $\mathfrak{g}(-r+k) \otimes \mathfrak{g}(r) = V_1 \oplus \dots \oplus V_m$ be the irreducible decomposition of $\mathfrak{g}(-r+k) \otimes \mathfrak{g}(r)$ as an L -module. Covariant map τ_k induces an L -equivariant linear map $\tilde{\tau}_k|_{V_j^*} : V_j^* \rightarrow \mathcal{P}^k(\mathfrak{g}(1))$ with V_j^* the dual of an irreducible constituent V_j of $\mathfrak{g}(-r+k) \otimes \mathfrak{g}(r)$ and $\mathcal{P}^k(\mathfrak{g}(1))$ the space of homogeneous polynomials on $\mathfrak{g}(1)$ of degree k . We define differential operators from $\tilde{\tau}_k|_{V_j^*}(Y^*)$. For $Y^* \in V_j^*$, let $\Omega_k(Y^*)$ denote the k -th order differential operators that are constructed from $\tilde{\tau}_k|_{V_j^*}(Y^*)$. We say that a list of differential operators D_1, \dots, D_n is the $\Omega_k|_{V_j^*}$ system if it is equivalent (see Definition 2.4) to a list of differential operators $\Omega_k(Y_1^*), \dots, \Omega_k(Y_n^*)$, where $\{Y_1^*, \dots, Y_n^*\}$ is a basis for V_j^* . By construction the $\Omega_k|_{V_j^*}$ system consists of $\dim_{\mathbb{C}}(V_j)$ operators.

The conformal invariance of the $\Omega_k|_{V_j^*}$ system depends on the complex parameter s for the line bundle \mathcal{L}_s . Then we say that the $\Omega_k|_{V_j^*}$ system has *special value* s_0 if the system is conformally invariant on the line bundle \mathcal{L}_{s_0} . The special values for the case that $\dim([\mathfrak{n}, \mathfrak{n}]) = 1$ for $\mathfrak{q} = \mathfrak{l} \oplus \mathfrak{n}$ are studied by Barchini-Kable-Zierau in [1] and [2], and myself in [17] and [18].

In this paper we consider a more general case; namely, $\mathfrak{q} = \mathfrak{l} \oplus \mathfrak{n}$ is a maximal parabolic subalgebra and \mathfrak{n} satisfies the condition that $[\mathfrak{n}, [\mathfrak{n}, \mathfrak{n}]] = 0$ and $\dim_{\mathbb{C}}([\mathfrak{n}, \mathfrak{n}]) > 1$. We call such parabolic subalgebras maximal two-step nilpotent parabolic subalgebras of non-Heisenberg type. In this case we have $r = 2$ in (3.4). Therefore the Ω_k systems for $k \geq 5$ are zero. We determine the special values of the Ω_1 system and Ω_2 systems associated to the parabolic subalgebras under consideration.

We may want to remark that, although the special value of s for the Ω_1 system is easily found by computing the bracket $[X, \Omega_1(Y_i^*)]$, it is not easy to find the special values for the Ω_2 systems by a direct computation. (See Section 5 of [1].) In this paper, to find the special values for the Ω_2 systems, we use two reduction techniques. These techniques significantly reduce the amount of computations. (See Proposition 2.16 and Proposition 7.13.)

This paper consists of seven sections including this introduction and one appendix. We now outline the contents of the rest of this paper. In Section 2, we first recall the definition of conformally invariant systems of differential operators and then collect some useful formulas. In Section 3, the construction of the Ω_k systems is given precisely. To construct the Ω_1 system and Ω_2 systems for maximal two-step nilpotent parabolic subalgebra \mathfrak{q} of non-Heisenberg type, we study such a parabolic subalgebra \mathfrak{q} and the associated 2-grading $\mathfrak{g} = \bigoplus_{j=-2}^2 \mathfrak{g}(j)$ in Section 4.

We construct the Ω_1 system and find its special value in Section 5. In this section we also fix normalizations for root vectors. The normalizations play an important role to construct the system. We show that the special value s_1 for the Ω_1 system is $s_1 = 0$. This is done in Theorem 5.7.

To build the Ω_2 systems, we need to find the irreducible constituents V^* of $\mathfrak{l}^* \otimes \mathfrak{z}(\mathfrak{n})^*$ so that $\tilde{\tau}_2|_{V^*} \neq 0$. In Section 6, we show preliminary results to find such irreducible constituents. First we decompose $\mathfrak{l} \otimes \mathfrak{z}(\mathfrak{n})$ into the direct sum of the irreducible constituents. By using the decomposition results, we then determine the candidates of the irreducible constituents V^* so that $\tilde{\tau}_2|_{V^*} \neq 0$. We build the Ω_2 systems and find their special values in Section 7. The special values are determined in Theorem 7.16.

Finally, in Appendix A, we summarize the miscellaneous useful data for the parabolic subalgebras under consideration.

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2. CONFORMALLY INVARIANT SYSTEMS

In this section we recall from [2] the definition of a conformally invariant system of differential operators. We also collect some properties of such a system of differential operators.

2.1. Conformally Invariant Systems. Let \mathfrak{g}_0 be a real Lie algebra and $\mathcal{X}(M)$ be the space of smooth vector fields on a smooth manifold M .

Definition 2.1. [2, page 790] *A smooth manifold M is called a \mathfrak{g}_0 -manifold if there is an \mathbb{R} -linear map $\pi_M : \mathfrak{g}_0 \rightarrow C^\infty(M) \oplus \mathcal{X}(M)$ so that for all $X, Y \in \mathfrak{g}_0$,*

$$\pi_M([X, Y]) = [\pi_M(X), \pi_M(Y)].$$

For each $X \in \mathfrak{g}_0$, we write $\pi_M(X) = \pi_0(X) + \pi_1(X)$ with $\pi_0(X) \in C^\infty(M)$ and $\pi_1(X) \in \mathcal{X}(M)$. We denote by $\mathbb{D}(\mathcal{V})$ the space of differential operators on smooth sections of \mathcal{V} .

Definition 2.2. [2, page 791] *Let M be a \mathfrak{g}_0 -manifold. A vector bundle $\mathcal{V} \rightarrow M$ is called a \mathfrak{g}_0 -bundle if there is an \mathbb{R} -linear map $\pi_{\mathcal{V}} : \mathfrak{g}_0 \rightarrow \mathbb{D}(\mathcal{V})$ that satisfies the following properties:*

- (B1) We have $\pi_{\mathcal{V}}([X, Y]) = [\pi_{\mathcal{V}}(X), \pi_{\mathcal{V}}(Y)]$ for all $X, Y \in \mathfrak{g}_0$.
 (B2) In $\mathbb{D}(\mathcal{V})$, $[\pi_{\mathcal{V}}(X), f] = \pi_1(X) \bullet f$ for all $X \in \mathfrak{g}_0$ and $f \in C^\infty(M)$.

Definition 2.3. [2, page 791] Let $\mathcal{V} \rightarrow M$ be a \mathfrak{g}_0 -bundle. A **conformally invariant system** on \mathcal{V} with respect to $\pi_{\mathcal{V}}$ is a list of differential operators $D_1, \dots, D_m \in \mathbb{D}(\mathcal{V})$ so that the following two conditions hold:

- (S1) At each point $p \in M$, the list D_1, \dots, D_m is linearly independent.
 (S2) For each $X \in \mathfrak{g}_0$, there is a matrix $C(X)$ in $M_{m \times m}(C^\infty(M))$ so that in $\mathbb{D}(\mathcal{V})$,

$$[\pi_{\mathcal{V}}(X), D_i] = \sum_{j=1}^m C_{ji}(X) D_j.$$

The map $C : \mathfrak{g}_0 \rightarrow M_{m \times m}(C^\infty(M))$ is called the **structure operator** of the conformally invariant system.

If \mathfrak{g} is the complexification of \mathfrak{g}_0 then \mathfrak{g} -manifolds and \mathfrak{g} -bundles are defined by extending the \mathfrak{g}_0 -action \mathbb{C} -linearly. In p. 792 in [2], the equivalence of two conformally invariant systems are defined. For later convenience we apply the same definition to any systems of differential operators. (See Definition 3.7.)

Definition 2.4. We say that two systems of differential operators (not necessarily conformally invariant) D_1, \dots, D_n and D'_1, \dots, D'_n in $\mathbb{D}(\mathcal{V})$ are **equivalent** if there is a matrix $A \in GL(n, \mathbb{C}^\infty(M))$ so that, for $1 \leq i \leq n$,

$$D'_i = \sum_{j=1}^n A_{ji} D_j.$$

Definition 2.5. [2, page 793] A conformally invariant system D_1, \dots, D_n is called **reducible** if there are an equivalent conformally invariant system D'_1, \dots, D'_n and an $m < n$ such that the system D'_1, \dots, D'_m is conformally invariant. Otherwise we say that D_1, \dots, D_n is **irreducible**.

We now specialize the \mathfrak{g} -manifold and \mathfrak{g} -bundle that we will work with. Let G be a complex, simple, connected, simply-connected Lie group with Lie algebra \mathfrak{g} . Such G contains a maximal connected solvable subgroup B . Write $\mathfrak{b} = \mathfrak{h} \oplus \mathfrak{u}$ for its Lie algebra with \mathfrak{h} the Cartan subalgebra and \mathfrak{u} the nilpotent subalgebra. Let $\mathfrak{q} \supset \mathfrak{b}$ be a parabolic subalgebra of \mathfrak{g} . We define $Q = N_G(\mathfrak{q})$, a parabolic subgroup of G . Write $Q = LN$ for the Levi decomposition of Q . with L the Levi subgroup and N the nilpotent subgroup.

Let \mathfrak{g}_0 be a real form of \mathfrak{g} and let G_0 be the analytic subgroup of G with Lie algebra \mathfrak{g}_0 . Define $Q_0 = N_{G_0}(\mathfrak{q}) \subset Q$, and write $Q_0 = L_0 N_0$. We will work with $M = G_0/Q_0$ for a class of maximal parabolic subgroup Q_0 that will be specified in Section 4.

Next, we need to specify a vector bundle \mathcal{V} on M . To this end let $\Delta = \Delta(\mathfrak{g}, \mathfrak{h})$ be the set of roots of \mathfrak{g} with respect to \mathfrak{h} . Let Δ^+ be the positive system attached to \mathfrak{b} and denote by Π the set of simple roots. For each subset $S \subset \Pi$, let \mathfrak{q}_S be the corresponding standard parabolic subalgebra. Write $\mathfrak{q}_S = \mathfrak{l}_S \oplus \mathfrak{n}_S$ with Levi factor $\mathfrak{l}_S = \mathfrak{h} \oplus \bigoplus_{\alpha \in \Delta_S} \mathfrak{g}_\alpha$ and nilpotent radical $\mathfrak{n}_S = \bigoplus_{\alpha \in \Delta^+ \setminus \Delta_S} \mathfrak{g}_\alpha$, where $\Delta_S = \{\alpha \in \Delta \mid \alpha \in \text{span}(\Pi \setminus S)\}$. If Q_0 is a maximal parabolic then there exists a unique simple root $\alpha_q \in \Pi$ so that $\mathfrak{q} = \mathfrak{q}_{\{\alpha_q\}}$. Let λ_q be the fundamental weight of α_q . The weight λ_q

is orthogonal to any roots α with $\mathfrak{g}_\alpha \subset [\mathfrak{l}, \mathfrak{l}]$. Hence it exponentiates to a character $\chi_{\mathfrak{q}}$ of L . As $\chi_{\mathfrak{q}}$ takes real values on L_0 , for $s \in \mathbb{C}$, character $\chi^s = |\chi_{\mathfrak{q}}|^s$ is well-defined on L_0 . Let \mathbb{C}_{χ^s} be the one-dimensional representation of L_0 with character χ^s . The representation χ^s is extended to a representation of Q_0 by making it trivial on N_0 . Then it deduces a line bundle \mathcal{L}_s on $M = G_0/Q_0$ with fiber \mathbb{C}_{χ^s} .

The group G_0 acts on the space

$$\begin{aligned} & C_\chi^\infty(G_0/Q_0, \mathbb{C}_{\chi^s}) \\ &= \{F \in C^\infty(G_0, \mathbb{C}_{\chi^s}) \mid F(gq) = \chi^s(q^{-1})F(g) \text{ for all } q \in Q_0 \text{ and } g \in G_0\} \end{aligned}$$

by left translation. The action π_s of \mathfrak{g}_0 on $C_\chi^\infty(G_0/Q_0, \mathbb{C}_{\chi^s})$ arising from this action is given by

$$(\pi_s(Y) \bullet F)(g) = \frac{d}{dt} F(\exp(-tY)g) \Big|_{t=0}$$

for $Y \in \mathfrak{g}_0$. This action is extended \mathbb{C} -linearly to \mathfrak{g} and then naturally to the universal enveloping algebra $\mathcal{U}(\mathfrak{g})$. We use the same symbols for the extended actions.

Let \bar{N}_0 be the nilpotent subgroup opposite to N_0 . By the Bruhat theory, the subset $\bar{N}_0 Q_0$ is open and dense in G_0 . Then the restriction map $C_\chi^\infty(G_0/Q_0, \mathbb{C}_{\chi^s}) \rightarrow C^\infty(\bar{N}_0, \mathbb{C}_{\chi^s})$ is an injection, where $C^\infty(\bar{N}_0, \mathbb{C}_{\chi^s})$ is the space of the smooth functions from \bar{N}_0 to \mathbb{C}_{χ^s} . Then, for $u \in \mathcal{U}(\mathfrak{g})$ and $F \in C_\chi^\infty(G_0/Q_0, \mathbb{C}_{\chi^s})$, we let $f = F|_{\bar{N}_0}$ and define the action of $\mathcal{U}(\mathfrak{g})$ on the image of the restriction map by

$$(2.6) \quad \pi_s(u) \bullet f = (\pi_s(u) \bullet F)|_{\bar{N}_0}.$$

The line bundle $\mathcal{L}_s \rightarrow G_0/Q_0$ restricted to \bar{N}_0 is the trivial bundle $\bar{N}_0 \times \mathbb{C}_{\chi^s} \rightarrow \bar{N}_0$. By slight abuse of notation, we refer to the trivial bundle over \bar{N}_0 as \mathcal{L}_s . Then in practice our manifold M will be $M = \bar{N}_0$ and our vector bundle will be the trivial bundle.

Now we show that, with the action π_s , the group \bar{N}_0 and the trivial bundle \mathcal{L}_s are a \mathfrak{g} -manifold and \mathfrak{g} -bundle, respectively. Let $\bar{\mathfrak{n}}$ and \mathfrak{q} be the complexifications of the Lie algebras of \bar{N}_0 and Q_0 , respectively; we have the direct sum $\mathfrak{g} = \bar{\mathfrak{n}} \oplus \mathfrak{q}$. For $Y \in \mathfrak{g}$, write $Y = Y_{\bar{\mathfrak{n}}} + Y_{\mathfrak{q}}$ for the decomposition of Y in this direct sum. Similarly, write the Bruhat decomposition of $g \in \bar{N}_0 Q_0$ as $g = \bar{\mathbf{n}}(g)\mathbf{q}(g)$ with $\bar{\mathbf{n}}(g) \in \bar{N}_0$ and $\mathbf{q}(g) \in Q_0$. For $Y \in \mathfrak{g}_0$, we have

$$(2.7) \quad Y_{\bar{\mathfrak{n}}} = \frac{d}{dt} \bar{\mathbf{n}}(\exp(tY)) \Big|_{t=0},$$

and a similar equality holds for $Y_{\mathfrak{q}}$. Define a right action R of $\mathcal{U}(\bar{\mathfrak{n}})$ on $C^\infty(\bar{N}_0, \mathbb{C}_{\chi^s})$ by

$$(2.8) \quad (R(X) \bullet f)(\bar{n}) = \frac{d}{dt} f(\bar{n} \exp(tX)) \Big|_{t=0}$$

for $X \in \bar{\mathfrak{n}}_0$ and $f \in C^\infty(\bar{N}_0, \mathbb{C}_{\chi^s})$. Observe that, by definition, the differential $d\chi$ of χ is $d\chi = \lambda_{\mathfrak{q}}$. A direct computation then shows that, for $Y \in \mathfrak{g}$ and f in the image of the restriction map $C_\chi^\infty(G_0/Q_0, \mathbb{C}_{\chi^s}) \rightarrow C^\infty(\bar{N}_0, \mathbb{C}_{\chi^s})$, we have

$$(2.9) \quad (\pi_s(Y) \bullet f)(\bar{n}) = s\lambda_{\mathfrak{q}}((\text{Ad}(\bar{n}^{-1})Y)_{\mathfrak{q}})f(\bar{n}) - (R((\text{Ad}(\bar{n}^{-1})Y)_{\bar{\mathfrak{n}}}) \bullet f)(\bar{n}).$$

Observe that (2.9) implies that the representation π_s extends to a representation of $\mathcal{U}(\mathfrak{g})$ on the whole space $C^\infty(\bar{N}_0, \mathbb{C}_{\chi^s})$. Moreover, it also shows that for all $Y \in \mathfrak{g}$, the linear map $\pi_s(Y)$ is in $C^\infty(\bar{N}_0) \oplus \mathcal{X}(\bar{N}_0)$. Therefore, with this linear map π_s , \bar{N}_0 is a \mathfrak{g} -manifold.

Next, we show that the linear map π_s gives \mathcal{L}_s the structure of a \mathfrak{g} -bundle. As π_s is a representation of \mathfrak{g} , the condition (B1) of Definition 2.2 is trivial. Thus it suffices to show that the condition (B2) holds. Since \mathcal{L}_s is the trivial bundle of \bar{N}_0 with fiber \mathbb{C}_{χ^s} , the space of smooth sections of \mathcal{L}_s is identified with $C^\infty(\bar{N}_0, \mathbb{C}_{\chi^s})$. The following proposition establishes the condition (B2) in our situation.

Proposition 2.10. *In $\mathbb{D}(\mathcal{L}_s)$, for $Y \in \mathfrak{g}$ and $f \in C^\infty(\bar{N}_0)$, we have*

$$([\pi_s(Y), f])(\bar{n}) = -(R((\text{Ad}(\bar{n}^{-1})Y)_{\bar{n}}) \bullet f)(\bar{n}).$$

Proof. This follows from the definition of $[\pi_s(Y), f]$ and formula (2.9). \square

2.2. Properties of Conformally Invariant Systems. In Section 3 we are going to construct systems of differential operators on \mathcal{L}_s . The systems of operators will satisfy several properties of conformally invariant systems. For convenience we collect those properties from [2] here.

We first define an action of L_0 on $\mathbb{D}(\mathcal{L}_s)$. As on p. 805 of [2], we define an action of L_0 on $C^\infty(\bar{N}_0, \mathbb{C}_{\chi^s})$ by

$$(l \cdot f)(\bar{n}) = \chi^s(l) f(l^{-1} \bar{n} l).$$

This action agrees with the action of L_0 by the left translation on the image of the restriction map $C^\infty(G_0/Q_0, \mathbb{C}_{\chi^s}) \rightarrow C^\infty(\bar{N}_0, \mathbb{C}_{\chi^s})$. In terms of this action we define an action of L_0 on $\mathbb{D}(\mathcal{L}_s)$ by

$$(2.11) \quad (l \cdot D) \bullet f = l \cdot (D \bullet (l^{-1} \cdot f)).$$

Definition 2.12. [2, page 806] *A conformally invariant system D_1, \dots, D_m on $\mathcal{L}_s \rightarrow \bar{N}_0$ is called **L_0 -stable** if there is a map $c : L_0 \rightarrow GL(n, C^\infty(\bar{N}_0))$ such that*

$$l \cdot D_i = \sum_{j=1}^m c(l)_{ji} D_j.$$

It is known that there exists a semisimple element $H_0 \in \mathfrak{l}$, so that $\text{ad}(H_0)$ has only integer eigenvalues on \mathfrak{g} with $\mathfrak{g}(1) \neq \{0\}$, $\mathfrak{l} = \mathfrak{g}(0)$, $\mathfrak{n} = \bigoplus_{j>0} \mathfrak{g}(j)$, and $\bar{\mathfrak{n}} = \bigoplus_{j>0} \mathfrak{g}(-j)$, where $\mathfrak{g}(j)$ is the j -eigenspace of $\text{ad}(H_0)$. (See for example [12, Section X.3])

Definition 2.13. [2, page 804] *A conformally invariant system D_1, \dots, D_m is called **homogeneous** if $C(H_0)$ is a scalar matrix, where C is the structure operator of the conformally invariant system. (See Definition 2.3.)*

Proposition 2.14. [2, Proposition 17] *Any irreducible conformally invariant system is homogeneous.*

Define

$$\mathbb{D}(\mathcal{L}_s)^{\bar{\mathfrak{n}}} = \{D \in \mathbb{D}(\mathcal{L}_s) \mid [\pi_s(X), D] = 0 \text{ for all } X \in \bar{\mathfrak{n}}\}.$$

Observe that in the sense of [2, page 796], the \mathfrak{g} -manifold \bar{N}_0 is *straight* with respect to the subalgebra $\bar{\mathfrak{n}}$ of \mathfrak{g} ([2, page 799]). Then we state the definition of *straight* conformally invariant systems specialized to the present situation. (For the general definition see p. 797 of [2].)

Definition 2.15. We say that a conformally invariant system D_1, \dots, D_m is **straight** if $D_j \in \mathbb{D}(\mathcal{L}_s)^{\bar{\mathfrak{n}}}$ for $j = 1, \dots, m$.

In general, to show that a given list D_1, \dots, D_m of differential operators on \bar{N}_0 is a conformally invariant system, we need check (S2) of Definition 2.3 at each point of \bar{N}_0 . Proposition 2.16 below shows that in the case D_1, \dots, D_m in $\mathbb{D}(\mathcal{L}_s)^{\bar{\mathfrak{n}}}$, it suffices to check the condition only at the identity e .

Proposition 2.16. [2, Proposition 13] Let D_1, \dots, D_m be a list of operators in $\mathbb{D}(\mathcal{L}_s)^{\bar{\mathfrak{n}}}$. Suppose that the list is linearly independent at e and that there is a map $b : \mathfrak{g} \rightarrow \mathfrak{gl}(m, \mathbb{C})$ such that

$$([\pi_s(Y), D_i] \bullet f)(e) = \sum_{j=1}^m b(Y)_{ji} (D_j \bullet f)(e)$$

for all $Y \in \mathfrak{g}$, $f \in C^\infty(\bar{N}_0, \mathbb{C}_{\chi^s})$, and $1 \leq i \leq m$. Then D_1, \dots, D_m is a conformally invariant system on \mathcal{L}_s . The structure operator of the system is given by $C(Y)(\bar{n}) = b(\text{Ad}(\bar{n}^{-1})Y)$ for all $\bar{n} \in \bar{N}_0$ and $Y \in \mathfrak{g}$.

2.3. Useful Formulas. To end this section we are going to show two formulas that will make certain arguments simple in Section 5 and Section 7.

Proposition 2.17. Let $Y \in \mathfrak{g}$ and $f \in C^\infty(\bar{N}_0, \mathbb{C}_{\chi^s})$. For $X, X_1, X_2 \in \bar{\mathfrak{n}}$, we have

$$([\pi_s(Y), R(X)] \bullet f)(\bar{n}) = (R([\text{Ad}(\bar{n}^{-1})Y]_{\mathfrak{q}}, X]_{\bar{\mathfrak{n}}}) \bullet f)(\bar{n}) - s\lambda_{\mathfrak{q}}([\text{Ad}(\bar{n}^{-1})Y, X]_{\mathfrak{q}})f(\bar{n})$$

and

$$\begin{aligned} & ([\pi_s(Y), R(X_1)R(X_2)] \bullet f)(\bar{n}) \\ &= (R([\text{Ad}(\bar{n}^{-1})Y]_{\mathfrak{q}}, X_1]_{\bar{\mathfrak{n}}})R(X_2) \bullet f)(\bar{n}) + (R(X_1)R([\text{Ad}(\bar{n}^{-1})Y]_{\mathfrak{q}}, X_2]_{\bar{\mathfrak{n}}}) \bullet f)(\bar{n}) \\ &+ (R([\text{Ad}(\bar{n}^{-1})Y, X_1]_{\mathfrak{q}}, X_2]_{\bar{\mathfrak{n}}}) \bullet f)(\bar{n}) - s\lambda_{\mathfrak{q}}([\text{Ad}(\bar{n}^{-1})Y, X_1]_{\mathfrak{q}})(R(X_2) \bullet f)(\bar{n}) \\ &- s\lambda_{\mathfrak{q}}([\text{Ad}(\bar{n}^{-1})Y, X_2]_{\mathfrak{q}})(R(X_1) \bullet f)(\bar{n}) - s\lambda_{\mathfrak{q}}([\text{Ad}(\bar{n}^{-1})Y, X_1], X_2]_{\mathfrak{q}})f(\bar{n}). \end{aligned}$$

Proof. These formulas follow from a direct computation on the left hand side of each equation using (2.9). \square

3. THE Ω_k SYSTEMS

The purpose of this section is to construct systems of k -th order differential operators in $\mathbb{D}(\mathcal{L}_s)^{\bar{\mathfrak{n}}}$ in a systematic manner. We shall call the systems of operators Ω_k systems.

3.1. Construction of the Ω_k systems. Let $\mathfrak{g} = \bigoplus_{j=-r}^r \mathfrak{g}(j)$ on \mathfrak{g} be a \mathbb{Z} -grading on \mathfrak{g} with $\mathfrak{g}(1) \neq 0$. By construction, $\mathfrak{q} = \mathfrak{g}(0) \oplus \bigoplus_{j>0} \mathfrak{g}(j)$ is a parabolic subalgebra. Take L to be the analytic subgroup of G with Lie algebra $\mathfrak{g}(0)$. Observe that, by Vinberg's Theorem ([12, Theorem 10.19]), the triple $(L, \text{Ad}, \mathfrak{g}(1))$ is a prehomogeneous vector space, that is, L has an open orbit in $\mathfrak{g}(1)$. To define our systems of differential operators, we use covariant maps (L -equivariant polynomial maps), which we denote by τ_k , associated to prehomogeneous vector space $(L, \text{Ad}, \mathfrak{g}(1))$. These maps can be thought to give symbols of a class of differential operators that we will study. We

would like to acknowledge that the construction of τ_k as in this paper was suggested by Anthony Kable.

Definition 3.1. Let $\mathfrak{g} = \bigoplus_{j=-r}^r \mathfrak{g}(j)$ be a graded complex simple Lie algebra with $\mathfrak{g}(1) \neq 0$. Then, for $0 \leq k \leq 2r$, the map τ_k on $\mathfrak{g}(1)$ is defined by

$$\begin{aligned} \tau_k : \mathfrak{g}(1) &\rightarrow \mathfrak{g}(-r+k) \otimes \mathfrak{g}(r) \\ X &\mapsto \frac{1}{k!} \text{ad}(X)^k \omega \end{aligned}$$

with $\omega = \sum_{\gamma_j \in \Delta(\mathfrak{g}(r))} X_{-\gamma_j} \otimes X_{\gamma_j}$, where X_{γ_j} are root vectors for γ_j and $\Delta(\mathfrak{g}(r))$ is the set of roots α so that $\mathfrak{g}_\alpha \subset \mathfrak{g}(r)$.

Here, we mean by $\text{ad}(X)^k \omega$ that X acts on the tensor product diagonally via the action $\text{ad}(\cdot)^k$. Observe that since $X \in \mathfrak{g}(1)$ and $[\mathfrak{g}(1), \mathfrak{g}(r)] = 0$, we have $\text{ad}(X)^k X_{\gamma_j} = 0$ for all $\gamma_j \in \Delta(\mathfrak{g}(r))$. Therefore, $\text{ad}(X)^k \omega = \sum_{\gamma_j} \text{ad}(X)^k (X_{-\gamma_j}) \otimes X_{\gamma_j}$.

We shall check in Lemma 3.3 that these maps are indeed L -equivariant. Observe that, by the standard argument, the element ω is independent of a choice of a basis for $\mathfrak{g}(r)$ and the dual basis for $\mathfrak{g}(-r)$.

Lemma 3.2. Let $\mathfrak{g} = \bigoplus_{j=-r}^r \mathfrak{g}(j)$ be a graded complex simple Lie algebra with $\mathfrak{g}(1) \neq 0$ and G be a complex analytic group with Lie algebra \mathfrak{g} . If L is the analytic subgroup of G with Lie algebra $\mathfrak{g}(0)$ and ω is as in Definition 3.1 then, for all $l \in L$,

$$(\text{Ad}(l) \otimes \text{Ad}(l))\omega = \omega.$$

Proof. If $g \in L$ then $\{\text{Ad}(l)X_{\gamma_j} \mid \gamma_j \in \Delta(\mathfrak{g}(r))\}$ forms a basis for $\mathfrak{g}(r)$. It also holds that $\{\text{Ad}(l)X_{-\gamma_j} \mid \gamma_j \in \Delta(\mathfrak{g}(r))\}$ is the dual basis for $\mathfrak{g}(-r)$ with respect to the Killing form. Now the assertion follows from the property that ω is independent of a choice of a basis for $\mathfrak{g}(r)$ and the dual basis for $\mathfrak{g}(-r)$. \square

Lemma 3.3. Let $\mathfrak{g} = \bigoplus_{j=-r}^r \mathfrak{g}(j)$, G , and L be as in Lemma 3.2. For all $l \in L$, $X \in \mathfrak{g}(1)$, and for $0 \leq k \leq 2r$, we have

$$\tau_k(\text{Ad}(l)X) = (\text{Ad}(l) \otimes \text{Ad}(l))\tau_k(X).$$

Proof. For $l \in L$, we have

$$\begin{aligned} \tau_k(\text{Ad}(l)X) &= \frac{1}{k!} \sum_{\gamma_j \in \Delta(\mathfrak{g}(r))} \text{ad}(\text{Ad}(l)(X))^k (X_{-\gamma_j}) \otimes X_{\gamma_j} \\ &= \frac{1}{k!} \sum_{\gamma_j \in \Delta(\mathfrak{g}(r))} \text{Ad}(l)(\text{ad}(X)^k (\text{Ad}(l^{-1})X_{-\gamma_j})) \otimes X_{\gamma_j} \\ &= (\text{Ad}(l) \otimes \text{Ad}(l)) \left(\frac{1}{k!} \sum_{\gamma_j \in \Delta(\mathfrak{g}(r))} \text{ad}(X)^k (\text{Ad}(l^{-1})X_{-\gamma_j}) \otimes \text{Ad}(l^{-1})(X_{\gamma_j}) \right) \\ &= (\text{Ad}(l) \otimes \text{Ad}(l)) \left(\frac{1}{k!} \text{ad}(X)^k \omega \right) \\ &= (\text{Ad}(l) \otimes \text{Ad}(l))\tau_k(X). \end{aligned}$$

Note that Lemma 3.2 is applied from line four to line five. \square

Now we are going to build the systems of differential operators in $\mathbb{D}(\mathcal{L}_s)^{\bar{n}}$ that we study. Observe that, as $\tau_k : \mathfrak{g}(1) \rightarrow \mathfrak{g}(-r+k) \otimes \mathfrak{g}(r) = W$ are L -equivariant polynomial maps of degree k , the maps τ_k can be thought of as elements in $(\mathcal{P}^k(\mathfrak{g}(1)) \otimes W)^L$, where $\mathcal{P}^k(\mathfrak{g}(1))$ denotes the space of homogeneous polynomials on $\mathfrak{g}(1)$ of degree k . Then the isomorphism $(\mathcal{P}^k(\mathfrak{g}(1)) \otimes W)^L \cong \text{Hom}_L(W^*, \mathcal{P}^k(\mathfrak{g}(1)))$ yields the L -intertwining operators $\tilde{\tau}_k$, that are given by

$$(3.4) \quad \tilde{\tau}_k(Y^*)(X) = Y^*(\tau_k(X)),$$

where W^* is the dual module of W with respect to the Killing form. For each $Y^* \in W^*$, we have $\tilde{\tau}_k(Y^*) \in \mathcal{P}^k(\mathfrak{g}(1)) \cong \text{Sym}^k(\mathfrak{g}(-1))$. We define differential operators in $\mathbb{D}(\mathcal{L}_s)^{\bar{n}}$ from $\tilde{\tau}_k(Y^*)$. This is done as follows. Let $\sigma : \text{Sym}^k(\mathfrak{g}(-1)) \rightarrow \mathcal{U}(\bar{n})$ be the symmetrization operator. Identify $\mathcal{U}(\bar{n})$ with $\mathbb{D}(\mathcal{L}_s)^{\bar{n}}$ by making \bar{n} act on $C^\infty(\bar{N}_0, \mathbb{C}_{\chi^s})$ via right differentiation R . Then we have a composition of linear maps

$$W^* \xrightarrow{\tilde{\tau}_k} \mathcal{P}^k(\mathfrak{g}(1)) \cong \text{Sym}^k(\mathfrak{g}(-1)) \xrightarrow{\sigma} \mathcal{U}(\bar{n}) \xrightarrow{R} \mathbb{D}(\mathcal{L}_s)^{\bar{n}}.$$

For $Y^* \in W^*$, we define a differential operator $\Omega_k(Y^*) \in \mathbb{D}(\mathcal{L}_s)^{\bar{n}}$ by

$$\Omega_k(Y^*) = R \circ \sigma \circ \tilde{\tau}_k(Y^*).$$

As we will work with irreducible systems we need to be a little more careful with our construction; in particular, we need to take an irreducible constituent of $\mathfrak{g}(-r+k)^* \otimes \mathfrak{g}(r)^*$. Let $\mathfrak{g}(-r+k) \otimes \mathfrak{g}(r) = V_1 \oplus \dots \oplus V_m$ be the irreducible decomposition of $\mathfrak{g}(-r+k) \otimes \mathfrak{g}(r)$ as an L -module, and let $\mathfrak{g}(-r+k)^* \otimes \mathfrak{g}(r)^* = V_1^* \oplus \dots \oplus V_m^*$ be the corresponding irreducible decomposition of $\mathfrak{g}(-r+k)^* \otimes \mathfrak{g}(r)^*$, where $\mathfrak{g}(j)^*$ are the dual spaces of $\mathfrak{g}(j)$ with respect to the Killing form. For each irreducible constituent V_j^* of $\mathfrak{g}(-r+k)^* \otimes \mathfrak{g}(r)^*$, there exists an L -intertwining operator $\tilde{\tau}_k|_{V_j^*} \in \text{Hom}_L(V_j^*, \mathcal{P}^k(\mathfrak{g}(1)))$ given as in (3.4). Then we define a linear operator $\Omega_k|_{V_j^*} : V_j^* \rightarrow \mathbb{D}(\mathcal{L}_s)^{\bar{n}}$ by

$$\Omega_k|_{V_j^*} = R \circ \sigma \circ \tilde{\tau}_k|_{V_j^*}.$$

Since, for $Y^* \in V_j^*$, we have $\Omega_k|_{V_j^*}(Y^*) = \Omega_k(Y^*)$ as a differential operator, we simply write $\Omega_k(Y^*)$ for the differential operator arising from $Y^* \in V_j^*$.

Definition 3.5. *If V^* is an irreducible constituent of $\mathfrak{g}(-r+k)^* \otimes \mathfrak{g}(r)^*$ so that $\tilde{\tau}_k|_{V^*} \neq 0$ then a list of differential operators $D_1, \dots, D_n \in \mathbb{D}(\mathcal{L}_s)^{\bar{n}}$ is called the $\Omega_k|_{V^*}$ **system** if it is equivalent (see Definition 2.4) to a list of differential operators*

$$(3.6) \quad \Omega_k(Y_1^*), \dots, \Omega_k(Y_n^*),$$

where $\{Y_1^*, \dots, Y_n^*\}$ is a basis for V^* .

We also simply refer each $\Omega_k|_{W^*}$ system to an Ω_k system. We want to remark that the construction of the Ω_k systems might require additional modification to secure the conformal invariance. See Section 6 in [1] and Section 3 in [18] for the modification for the Ω_3 systems of the Heisenberg parabolic subalgebra.

It is important to notice that it is not necessary for the Ω_k systems to be conformally invariant; their conformal invariance strongly depends on the complex parameter s for the line bundle \mathcal{L}_s . So, we give the following definition.

Definition 3.7. Let V^* be an irreducible constituent of $\mathfrak{g}(-r+k)^* \otimes \mathfrak{g}(r)^*$. Then we say that the $\Omega_k|_{V^*}$ system has **special value** s_0 if the system is conformally invariant on the line bundle \mathcal{L}_{s_0} .

Note that, as the opposite parabolic $\tilde{Q}_0 = L_0\tilde{N}_0$ is chosen in [1], our special values s_0 are of the form $s_0 = -s'_0$, where s'_0 are the special values shown in [1].

Observe that the linear operator $\Omega_k|_{V^*} : V^* \rightarrow \mathbb{D}(\mathcal{L}_s)^{\bar{n}}$ is an L_0 -intertwining operator with respect to the action given in (2.11); in particular, the $\Omega_k|_{V^*}$ system is L_0 -stable (see Definition 2.12). Indeed, one can check that we have $l \cdot R(u) = R(\text{Ad}(l)u)$ for $l \in L_0$ and $u \in \mathcal{U}(\bar{n})$. This action stabilizes the subspace $\mathbb{D}(\mathcal{L}_s)^{\bar{n}}$. With the adjoint action of L_0 on $\mathcal{U}(\bar{n})$, the linear isomorphism $\mathcal{U}(\bar{n}) \xrightarrow{R} \mathbb{D}(\mathcal{L}_s)^{\bar{n}}$ is then L_0 -equivariant. It is clear that each map in $V^* \xrightarrow{\tilde{\tau}_k|_{Y^*}} \mathcal{P}^k(\mathfrak{g}(1)) \cong \text{Sym}^k(\mathfrak{g}(-1)) \xrightarrow{\sigma} \mathcal{U}(\bar{n})$ is L_0 -equivariant with respect to the natural actions of L_0 on each space, which are induced by the adjoint action of L_0 on \mathfrak{g} . Therefore, with the L_0 -action (2.11), the operator $\Omega_k|_{V^*} : V^* \rightarrow \mathbb{D}(\mathcal{L}_s)^{\bar{n}}$ is an L_0 -intertwining operator. Now we summarize some properties of the $\Omega_k|_{V^*}$ system.

Remark 3.8. It follows from the definition and observation above that the $\Omega_k|_{V^*}$ system satisfies the following properties:

- (1) The $\Omega_k|_{V^*}$ system satisfies the condition (S1) of Definition 2.3.
- (2) When the $\Omega_k|_{V^*}$ system is conformally invariant then it is an irreducible, straight, and L_0 -stable system. By Proposition 2.14, it is also a homogeneous system.

3.2. Computations involving the Ω_k systems. We are going to show two technical lemmas that will be used in Section 7. For $D \in \mathbb{D}(\mathcal{L}_s)$, we denote by $D_{\bar{n}}$ the linear functional $f \mapsto (D \bullet f)(\bar{n})$ for $f \in C^\infty(\bar{N}_0, \mathbb{C}_{\chi^s})$. A simple observation shows that $(D_1 D_2)_{\bar{n}} = (D_1)_{\bar{n}} D_2$ for $D_1, D_2 \in \mathbb{D}(\mathcal{L}_s)$; in particular, if $(D_1)_{\bar{n}} = 0$ then $[D_1, D_2]_{\bar{n}} = -(D_2)_{\bar{n}} D_1$.

Lemma 3.9. Suppose that V^* is an irreducible constituent of $\mathfrak{g}(-r+k)^* \otimes \mathfrak{g}(r)^*$. Let $X_1, X_2 \in \mathfrak{g}$ and $Y_1^*, \dots, Y_n^* \in V^*$. If $\pi_s(X_1)_e = 0$ and if we have $[\pi_s(X_i), \Omega_k(Y_t^*)]_e \in \text{span}_{\mathbb{C}}\{\Omega_k(Y_j^*)_e \mid j = 1, \dots, n\}$ for $i = 1, 2$ then

$$(3.10) \quad [\pi_s(X_1), [\pi_s(X_2), \Omega_k(Y_t^*)]]_e \in \text{span}_{\mathbb{C}}\{\Omega_k(Y_1^*)_e, \dots, \Omega_k(Y_n^*)_e\}.$$

Proof. Observe that $[\pi_s(X_1), [\pi_s(X_2), \Omega_k(Y_t^*)]]$ is

$$(3.11) \quad \pi_s(X_1)[\pi_s(X_2), \Omega_k(Y_t^*)] - [\pi_s(X_2), \Omega_k(Y_t^*)]\pi_s(X_1).$$

Since, by assumption, we have $\pi_s(X_1)_e = 0$, the first term is zero at e . By assumption, the bracket $[\pi_s(X_2), \Omega_k(Y_t^*)]_e$ is a linear combination of $\Omega_k(Y_1^*)_e, \dots, \Omega_k(Y_n^*)_e$ over \mathbb{C} . So it may be written as $[\pi_s(X_2), \Omega_k(Y_t^*)]_e = \sum_{j=1}^n a_{jt} \Omega_k(Y_j^*)_e$ with $a_{jt} \in \mathbb{C}$. Then, the second term in (3.11) evaluates to $-\sum_{j=1}^n a_{jt} \Omega_k(Y_j^*)_e \pi_s(X_1)$ at the identity e . Since $(\pi_s(X_1) \Omega_k(Y_j^*))_e = \pi_s(X_1)_e \Omega_k(Y_j^*)_e = 0$, we

obtain

$$\begin{aligned} [\pi_s(X_1), [\pi_s(X_2), \Omega_k(Y_t^*)]]_e &= - \sum_{j=1}^n a_{jt} \Omega_k(Y_j^*)_e \pi_s(X_1) \\ &= - \sum_{j=1}^n a_{jt} [\pi_s(X_1), \Omega_k(Y_j^*)]_e. \end{aligned}$$

Now the proposed result follows from the assumption that $[\pi_s(X_1), \Omega_k(Y_t^*)]_e$ is a linear combination of $\Omega_k(Y_j^*)_e$ over \mathbb{C} . \square

If $\Delta^+(\mathfrak{l})$ is the set of positive roots in \mathfrak{l} then we set

$$\mathfrak{u}_{\mathfrak{l}} = \bigoplus_{\Delta^+(\mathfrak{l})} \mathfrak{g}_{\alpha} \text{ and } \bar{\mathfrak{u}}_{\mathfrak{l}} = \bigoplus_{\Delta^+(\mathfrak{l})} \mathfrak{g}_{-\alpha}.$$

Lemma 3.12. *Suppose that $\mathfrak{g}(1)$ is irreducible and that V^* is an irreducible constituent of $\mathfrak{g}(-r + k)^* \otimes \mathfrak{g}(r)^*$. Let X_h be a highest weight vector for $\mathfrak{g}(1)$ and Y_l^* be a lowest weight vector for V^* . If*

$$[\pi_s(X_h), \Omega_k(Y_l^*)]_e \in \text{span}_{\mathbb{C}}\{\Omega_k(Y_1^*)_e, \dots, \Omega_k(Y_n^*)_e\}$$

with $\{Y_1^, \dots, Y_n^*\}$ a basis for V^* then, for any $X \in \mathfrak{g}(1)$ and $Y^* \in V^*$,*

$$[\pi_s(X), \Omega_k(Y^*)]_e \in \text{span}_{\mathbb{C}}\{\Omega_k(Y_1^*)_e, \dots, \Omega_k(Y_n^*)_e\}.$$

Proof. Set $E = \text{span}_{\mathbb{C}}\{\Omega_k(Y_1^*)_e, \dots, \Omega_k(Y_n^*)_e\}$. We first show that for each $X \in \mathfrak{g}(1)$,

$$(3.13) \quad [\pi_s(X), \Omega_k(Y_l^*)]_e \in E.$$

Observe that since $(L, \mathfrak{g}(1))$ is assumed to be irreducible, the L -module $\mathfrak{g}(1)$ is given by $\mathfrak{g}(1) = \mathcal{U}(\bar{\mathfrak{u}}_{\mathfrak{l}})X_h$. Then, as π_s is linear on $\mathfrak{g}(1)$, it suffices to show that (3.13) holds when $X = \bar{u}_k \cdot X_h$ with \bar{u}_k a monomial in $\mathcal{U}(\bar{\mathfrak{u}}_{\mathfrak{l}})$. This is done by induction on the order of \bar{u}_k . Indeed, the proof is clear once we show that (3.13) holds for $X = \bar{Z} \cdot X_h = [\bar{Z}, X_h]$ with $\bar{Z} \in \bar{\mathfrak{u}}_{\mathfrak{l}}$.

By the Jacobi identity, the commutator $[\pi_s([\bar{Z}, X_h]), \Omega_k(Y_l^*)]$ is

$$(3.14) \quad \begin{aligned} &[\pi_s([\bar{Z}, X_h]), \Omega_k(Y_l^*)] \\ &= [\pi_s(\bar{Z}), [\pi_s(X_h), \Omega_k(Y_l^*)]] - [\pi_s(X_h), [\pi_s(\bar{Z}), \Omega_k(Y_l^*)]]. \end{aligned}$$

By the \mathfrak{l} -equivariance of the operator $\Omega_k : V^* \rightarrow \mathbb{D}(\mathcal{L}_s)^{\bar{\mathfrak{n}}}$, it follows that

$$[\pi_s(\bar{Z}), \Omega_k(Y_l^*)] = \Omega_k([\bar{Z}, Y_l^*]).$$

Since $\bar{Z} \in \bar{\mathfrak{u}}_{\mathfrak{l}}$ and Y_l^* is a lowest weight vector, we have $\Omega_k([\bar{Z}, Y_l^*]) = 0$, and so is the second term of the right hand side of (3.14). Thus we have

$$(3.15) \quad [\pi_s([\bar{Z}, X_h]), \Omega_k(Y_l^*)]_e = [\pi_s(\bar{Z}), [\pi_s(X_h), \Omega_k(Y_l^*)]]_e.$$

Now, by hypotheses and the \mathfrak{l} -equivariance of Ω_k , it follows that

$$[\pi_s(X_h), \Omega_k(Y_l^*)]_e, [\pi_s(\bar{Z}), \Omega_k(Y_l^*)]_e \in E.$$

As $\bar{Z} \in \bar{\mathfrak{u}}_{\mathfrak{l}}$, by (2.9), we have $\pi_s(\bar{Z})_e = 0$. Thus it follows from Lemma 3.9 that we have $[\pi_s(\bar{Z}), [\pi_s(X_h), \Omega_k(Y_l^*)]]_e \in E$. Therefore, $[\pi_s([\bar{Z}, X_h]), \Omega_k(Y_l^*)]_e \in E$ by (3.15).

Next we show that for any $X \in \mathfrak{g}(1)$ and $Y^* \in V^*$,

$$(3.16) \quad [\pi_s(X), \Omega_k(Y^*)]_e \in E.$$

Once again since V^* is irreducible, it is given by $V^* = \mathcal{U}(\mathfrak{u}_l)Y_l^*$. As before, it is enough to show that (3.16) holds for $Y^* = Z \cdot Y_l^*$ with $Z \in \mathfrak{u}_l$. Since $\Omega_k(Z \cdot Y_l^*) = [\pi_s(Z), \Omega_k(Y_l^*)]$, by the Jacobi identity, the commutator $[\pi_s(X), \Omega_k(Z \cdot Y_l^*)]$ is

$$(3.17) \quad \begin{aligned} & [\pi_s(X), \Omega_k(Z \cdot Y_l^*)] \\ &= [\pi_s(Z), [\pi_s(X), \Omega_k(Y_l^*)]] - [[\pi_s(Z), \pi_s(X)], \Omega_k(Y_l^*)]. \end{aligned}$$

We showed above that $[\pi_s(X), \Omega_k(Y_l^*)]_e \in E$. Since we have $\pi_s(Z)_e = 0$ and $[\pi_s(Z), \Omega_k(Y_l^*)]_e \in E$, by Lemma 3.9, the first term of the right hand side of (3.17) satisfies

$$[\pi_s(Z), [\pi_s(X), \Omega_k(Y_l^*)]]_e \in E.$$

Moreover, as $[\pi_s(Z), \pi_s(X)] = \pi_s([Z, X])$ with $[Z, X] \in \mathfrak{g}(1)$, by what we have shown above, the second term satisfies

$$[[\pi_s(Z), \pi_s(X)], \Omega_k(Y_l^*)]_e \in E.$$

Hence, $[\pi_s(X), \Omega_k(Z \cdot Y_l^*)]_e \in E$. □

4. PARABOLIC SUBALGEBRAS AND \mathbb{Z} -GRADINGS

It has been observed in Subsection 3.1 that the \mathbb{Z} -grading $\mathfrak{g} = \bigoplus_{j=-r}^r \mathfrak{g}(j)$ on \mathfrak{g} and parabolic subalgebra \mathfrak{q} play a role to construct the Ω_k systems. In this section we observe these in detail for \mathfrak{q} a maximal two-step nilpotent parabolic subalgebra of non-Heisenberg type. The Ω_1 system and Ω_2 systems of those parabolic subalgebras will be constructed in Section 5 and Section 7, respectively.

4.1. k -step Nilpotent Parabolic Subalgebras. Let \mathfrak{r} be any nonzero Lie algebra. Put $\mathfrak{r}_0 = \mathfrak{r}$, $\mathfrak{r}_1 = [\mathfrak{r}, \mathfrak{r}]$, and $\mathfrak{r}_k = [\mathfrak{r}, \mathfrak{r}_{k-1}]$ for $k \in \mathbb{Z}_{>0}$. We call \mathfrak{r}_k the k -th step of \mathfrak{r} for $k \in \mathbb{Z}_{\geq 0}$. The Lie algebra \mathfrak{r} is called **nilpotent** if $\mathfrak{r}_k = 0$ for some k , and it is called **k -step nilpotent** if $\mathfrak{r}_{k-1} \neq 0$ and $\mathfrak{r}_k = 0$. In particular, if $[\mathfrak{r}, \mathfrak{r}] = 0$ then \mathfrak{r} is called **abelian**, and if $\dim([\mathfrak{r}, \mathfrak{r}]) = 1$ then \mathfrak{r} is called **Heisenberg**. Note that \mathfrak{r} is Heisenberg if and only if its center $\mathfrak{z}(\mathfrak{r})$ is one-dimensional. If the nilpotent radical \mathfrak{n} of a parabolic subalgebra $\mathfrak{q} = \mathfrak{l} \oplus \mathfrak{n}$ is k -step nilpotent (resp. abelian or Heisenberg) then we say that \mathfrak{q} is a **k -step nilpotent** (resp. **abelian** or **Heisenberg**) parabolic subalgebra.

To build the Ω_1 system and Ω_2 systems of a maximal two-step nilpotent parabolic subalgebra \mathfrak{q} , it is convenient to classify the k -step nilpotent parabolic subalgebras \mathfrak{q} by the subsets of simple roots. If $\beta = \sum_{\alpha \in \Pi} m_\alpha \alpha \in \sum_{\alpha \in \Pi} \mathbb{Z}\alpha$ then we say that $|m_\alpha|$ are the *multiplicities* of α in β .

Proposition 4.1. *Let \mathfrak{g} be a complex simple Lie algebra with highest root γ , and $\mathfrak{q}_S = \mathfrak{l} \oplus \mathfrak{n}$ be the parabolic subalgebra of \mathfrak{g} that is parametrized by S with $S = \{\alpha_{i_1}, \dots, \alpha_{i_r}\} \subset \Pi$. Then \mathfrak{n} is k -step nilpotent if and only if $k = m_{i_1} + m_{i_2} + \dots + m_{i_r}$, where m_{i_j} are the multiplicities of α_{i_j} in γ .*

Proof. As this is a well-known fact, we omit a proof. (For a proof, see for instance Section 3.1 of [15].) □

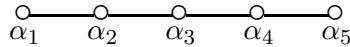
The following observation would be useful when we consider k -step nilpotent parabolic subalgebras. First, observe that, by the one-to-one correspondence between the standard parabolic subalgebras \mathfrak{q}_S and the subsets $S \subset \Pi$, we can associate subdiagrams of Dynkin diagrams to parabolic subalgebras \mathfrak{q}_S . The subdiagrams that associates to \mathfrak{q}_S are obtained by deleting the nodes of the Dynkin diagram of \mathfrak{g} that correspond to the simple roots in S , and the edges in incident on them. We call such subdiagrams **deleted Dynkin diagrams**. With the multiplicities of simple roots in the highest root of \mathfrak{g} in hand, by Proposition 4.1, we then see the number of steps of nilradical \mathfrak{n} of \mathfrak{q}_S from the deleted Dynkin diagram. Table 1 shows the multiplicities the simple roots in the highest root γ . We use the Bourbaki conventions [3] for the labels of the simple roots.

TABLE 1. Highest Roots

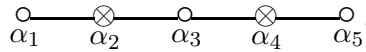
Type	Highest root
A_n	$\alpha_1 + \alpha_2 + \cdots + \alpha_n$
B_n	$\alpha_1 + 2\alpha_2 + 2\alpha_3 + \cdots + 2\alpha_n$
C_n	$2\alpha_1 + 2\alpha_2 + \cdots + 2\alpha_{n-1} + 1\alpha_n$
D_n	$\alpha_1 + 2\alpha_2 + 2\alpha_3 + \cdots + 2\alpha_{n-2} + \alpha_{n-1} + \alpha_n$
E_6	$\alpha_1 + 2\alpha_2 + 2\alpha_3 + 3\alpha_4 + 2\alpha_5 + \alpha_6$
E_7	$2\alpha_1 + 2\alpha_2 + 3\alpha_3 + 4\alpha_4 + 3\alpha_5 + 2\alpha_6 + \alpha_7$
E_8	$2\alpha_1 + 3\alpha_2 + 4\alpha_3 + 6\alpha_4 + 5\alpha_5 + 4\alpha_6 + 3\alpha_7 + 2\alpha_8$
F_4	$2\alpha_1 + 3\alpha_2 + 4\alpha_3 + 2\alpha_4$
G_2	$3\alpha_1 + 2\alpha_2$

Example 4.2 below describes the deleted Dynkin diagram of a given parabolic \mathfrak{q}_S and how we read the diagram. For simplicity, we depict deleted Dynkin diagrams by crossing out the deleted nodes.

Example 4.2. Let $\mathfrak{g} = \mathfrak{sl}(6, \mathbb{C})$. The Dynkin diagram is



Choose $S = \{\alpha_2, \alpha_4\} \subset \Pi$. Then the deleted Dynkin diagram of parabolic subalgebra \mathfrak{q}_S corresponding to the subset S is



Moreover, by Table 1, the multiplicity of each simple root in the highest root of \mathfrak{g} is 1. Thus, \mathfrak{q}_S is a two-step nilpotent parabolic subalgebra.

By the above observation we often refer to parabolic subalgebras \mathfrak{q}_S by their corresponding subset S of simple roots. To this end, we are going to define classification types of parabolics \mathfrak{q}_S . In Definition 4.3 below, we mean by classification type \mathcal{T} of \mathfrak{g} type A_n , B_n , C_n , D_n , E_6 , E_7 , E_8 , F_4 , or G_2 .

Definition 4.3. If \mathfrak{g} is a complex simple Lie algebra of classification type \mathcal{T} and if S is a subset of Π of simple roots then we say that a parabolic subalgebra \mathfrak{q}_S of \mathfrak{g} is of **type** $\mathcal{T}(S)$, or **type** $\mathcal{T}(i_1, \dots, i_k)$ if $S = \{\alpha_{i_1}, \dots, \alpha_{i_k}\}$.

For example, the parabolic subalgebra \mathfrak{q}_S in Example 4.2 is of type $A_5(2, 4)$. Any maximal parabolic subalgebra is of type $\mathcal{T}(i)$ for some $\alpha_i \in \Pi$.

4.2. Maximal Two-Step Nilpotent Parabolic \mathfrak{q} of Non-Heisenberg type. Now we observe the 2-grading $\mathfrak{g} = \bigoplus_{j=-2}^2 \mathfrak{g}(j)$ on \mathfrak{g} , that is induced from a maximal two-step nilpotent parabolic subalgebra \mathfrak{q} of non-Heisenberg type.

Assume that \mathfrak{g} has rank greater than one and that $\alpha_{\mathfrak{q}}$ is a simple root, so that the parabolic subalgebra $\mathfrak{q} = \mathfrak{q}_{\{\alpha_{\mathfrak{q}}\}} = \mathfrak{l} \oplus \mathfrak{n}$ parameterized by $\alpha_{\mathfrak{q}}$ is a maximal two-step nilpotent parabolic with $\dim([\mathfrak{n}, \mathfrak{n}]) > 1$. Let $\langle \cdot, \cdot \rangle$ be the inner product induced on \mathfrak{h}^* corresponding to the Killing form κ . Write $\|\alpha\|^2 = \langle \alpha, \alpha \rangle$ for $\alpha \in \Delta$. The coroot of α is $\alpha^\vee = 2\alpha/\langle \alpha, \alpha \rangle$.

Recall from Section 2 that $\lambda_{\mathfrak{q}}$ denotes the fundamental weight for $\alpha_{\mathfrak{q}}$. As $\Delta(\mathfrak{l}) = \{\alpha \in \Delta \mid \alpha \in \text{span}(\Pi \setminus \{\alpha_{\mathfrak{q}}\})\}$ and $\Delta(\mathfrak{n}) = \Delta^+ \setminus \Delta(\mathfrak{l})$, we have

$$\langle \lambda_{\mathfrak{q}}, \beta \rangle \begin{cases} = 0 & \text{if } \beta \in \Delta(\mathfrak{l}) \\ > 0 & \text{if } \beta \in \Delta(\mathfrak{n}) . \end{cases}$$

Observe that if $H_{\lambda_{\mathfrak{q}}} \in \mathfrak{h}$ is defined by $\kappa(H, H_{\lambda_{\mathfrak{q}}}) = \lambda_{\mathfrak{q}}(H)$ for all $H \in \mathfrak{h}$ and if

$$(4.4) \quad H_{\mathfrak{q}} = \frac{2}{\|\alpha_{\mathfrak{q}}\|^2} H_{\lambda_{\mathfrak{q}}}$$

then $\beta(H_{\mathfrak{q}})$ is the multiplicity of $\alpha_{\mathfrak{q}}$ in β . In particular, it follows from Proposition 4.1 that for $\beta \in \Delta^+$, $\beta(H_{\mathfrak{q}})$ can only take the values of 0, 1, or 2. Therefore, if $\mathfrak{g}(j)$ denotes the j -eigenspace of $\text{ad}(H_{\mathfrak{q}})$ then the action of $\text{ad}(H_{\mathfrak{q}})$ on \mathfrak{g} induces a 2-grading

$$\mathfrak{g} = \mathfrak{g}(-2) \oplus \mathfrak{g}(-1) \oplus \mathfrak{g}(0) \oplus \mathfrak{g}(1) \oplus \mathfrak{g}(2)$$

with parabolic subalgebra

$$\mathfrak{q} = \mathfrak{g}(0) \oplus \mathfrak{g}(1) \oplus \mathfrak{g}(2).$$

Here we have $\mathfrak{l} = \mathfrak{g}(0)$ and $\mathfrak{n} = \mathfrak{g}(1) \oplus \mathfrak{g}(2)$. The subalgebra $\bar{\mathfrak{n}}$, the opposite of \mathfrak{n} , is given by

$$\bar{\mathfrak{n}} = \mathfrak{g}(-1) \oplus \mathfrak{g}(-2).$$

Let $\mathfrak{l} = \mathfrak{z}(\mathfrak{l}) \oplus [\mathfrak{l}, \mathfrak{l}]$ be the decomposition of \mathfrak{l} , that corresponds to $L = Z(L)^\circ L_{ss}$ with $Z(L)^\circ$ the identity component of the center of L and L_{ss} the semisimple part of L . We say that a weight $\nu \in \mathfrak{h}^*$ is a highest weight of a finite dimensional L -module V if $\nu|_{\mathfrak{h}_{ss}}$ is a highest weight of V as an L_{ss} -module, where $\mathfrak{h}_{ss} = \mathfrak{h} \cap [\mathfrak{l}, \mathfrak{l}]$. A lowest weight of a finite dimensional L -module is similarly defined.

Proposition 4.5. *Let $\mathfrak{q} = \mathfrak{g}(0) \oplus \mathfrak{g}(1) \oplus \mathfrak{g}(2)$ be the maximal two-step nilpotent parabolic of non-Heisenberg type determined by $\alpha_{\mathfrak{q}}$.*

- (1) *The subspace $\mathfrak{g}(1)$ is the irreducible L -module with lowest weight $\alpha_{\mathfrak{q}}$.*
- (2) *The subspace $\mathfrak{g}(2)$ is the irreducible L -module with highest weight γ .*
- (3) *We have $\mathfrak{z}(\mathfrak{n}) = \mathfrak{g}(2)$.*

Proof. Observe that, as $\text{Ad}(L)$ preserves $\mathfrak{g}(j)$, to prove the assertions (1) and (2), it suffices to consider $\mathfrak{g}(1)$ and $\mathfrak{g}(2)$ as \mathfrak{l} -modules. For the assertion (1), the \mathfrak{l} -irreducibility of $\mathfrak{g}(1)$ just follows from a well-known fact that, for $\mathfrak{q} = \mathfrak{g}(0) \oplus \bigoplus_{j>0} \mathfrak{g}(j)$ with $\mathfrak{g}(1) \neq 0$, $\mathfrak{g}(1)$ is $\mathfrak{g}(0)$ -irreducible if and only if \mathfrak{q} is a maximal parabolic subalgebra. The lowest weight of $\mathfrak{g}(1)$ follows from Corollary 10.2A of [8]. For the assertion (2), it is clear that $\mathcal{U}(\mathfrak{g}(0))X_\gamma \subset \mathfrak{g}(2)$. On the other hand, as

$\bar{\mathfrak{n}} = \mathfrak{g}(-1) \oplus \mathfrak{g}(-2)$, it follows that $\mathcal{U}(\bar{\mathfrak{n}})\mathfrak{g}(2) \subset \bigoplus_{j=-2}^1 \mathfrak{g}(j)$. As $\mathfrak{g} = \bigoplus_{j=-2}^2 \mathfrak{g}(j) = \mathcal{U}(\bar{\mathfrak{n}})(\mathcal{U}(\mathfrak{g}(0))X_\gamma)$, this shows that $\mathcal{U}(\mathfrak{g}(0))X_\gamma \supset \mathfrak{g}(2)$. To prove assertion (3), since $\mathfrak{g}(2) \subset \mathfrak{z}(\mathfrak{n})$, it suffices to show the other inclusion. If $X \in \mathfrak{z}(\mathfrak{n})$ then, as $\mathfrak{n} = \mathfrak{g}(1) \oplus \mathfrak{g}(2)$, there exist $X_j \in \mathfrak{g}(j)$ for $j = 1, 2$ so that $X = X_1 + X_2$. Since $X, X_2 \in \mathfrak{z}(\mathfrak{n})$, for any $Y \in \mathfrak{n}$, we have

$$[Y, X_1] = [Y, X_1] + [Y, X_2] = [Y, X] = 0.$$

Thus $X_1 \in \mathfrak{z}(\mathfrak{n}) \cap \mathfrak{g}(1)$. Now observe that the assertion (1) implies that $\mathfrak{z}(\mathfrak{n}) \cap \mathfrak{g}(1) = \{0\}$. Thus $X_1 = 0$ and so $X = X_2 \in \mathfrak{g}(2)$. \square

Now, since $\mathfrak{l} = \mathfrak{g}(0)$, $\mathfrak{g}(2) = \mathfrak{z}(\mathfrak{n})$ and $\mathfrak{g}(-2) = \mathfrak{z}(\bar{\mathfrak{n}})$, we write the 2-grading $\mathfrak{g} = \bigoplus_{j=-2}^2 \mathfrak{g}(j)$ as

$$(4.6) \quad \mathfrak{g} = \mathfrak{z}(\bar{\mathfrak{n}}) \oplus \mathfrak{g}(-1) \oplus \mathfrak{l} \oplus \mathfrak{g}(1) \oplus \mathfrak{z}(\mathfrak{n})$$

with parabolic subalgebra

$$(4.7) \quad \mathfrak{q} = \mathfrak{l} \oplus \mathfrak{g}(1) \oplus \mathfrak{z}(\mathfrak{n}).$$

4.3. The simple ideals \mathfrak{l}_γ and $\mathfrak{l}_{n\gamma}$. We next observe the structure of the Levi subalgebra $\mathfrak{l} = \mathfrak{z}(\mathfrak{l}) \oplus [\mathfrak{l}, \mathfrak{l}]$. The structure of \mathfrak{l} will play a role in Section 6, when we decompose $\mathfrak{l} \otimes \mathfrak{z}(\mathfrak{n})$ into irreducible L -submodules.

We start with the center $\mathfrak{z}(\mathfrak{l})$. The center $\mathfrak{z}(\mathfrak{l})$ is of the form $\mathfrak{z}(\mathfrak{l}) = \bigcap_{\alpha \in \Pi(\mathfrak{l})} \ker(\alpha)$. Since \mathfrak{g} has rank greater than one and since $\Pi(\mathfrak{l}) = \Pi \setminus \{\alpha_q\}$, the center $\mathfrak{z}(\mathfrak{l})$ is non-zero and one-dimensional. It is clear from (4.4) that H_q is an element of $\mathfrak{z}(\mathfrak{l})$. Therefore we have $\mathfrak{z}(\mathfrak{l}) = \mathbb{C}H_q$.

Next we consider the structure of $[\mathfrak{l}, \mathfrak{l}]$. Observe that if \mathfrak{g} is not of type A_n then there is exactly one simple root that is not orthogonal to γ . Let α_γ denote the unique simple root. It is easy to see that $\mathfrak{q}_{\{\alpha_\gamma\}}$ is the Heisenberg parabolic of \mathfrak{g} ; that is, the two-step nilpotent parabolic subalgebra with $\dim([\mathfrak{n}, \mathfrak{n}]) = 1$. Hence, if $\mathfrak{q}_{\{\alpha_q\}}$ is a maximal two-step nilpotent parabolic subalgebra with $\dim([\mathfrak{n}, \mathfrak{n}]) > 1$ then $\alpha_\gamma \in \Pi(\mathfrak{l}) = \Pi \setminus \{\alpha_q\}$. If we delete the node corresponding to α_q then we obtain one, two, or three subgraphs with one subgraph containing α_γ . This implies that the subalgebra $[\mathfrak{l}, \mathfrak{l}]$ is either simple or the direct sum of two or three simple ideals with only one simple ideal containing the root space $\mathfrak{g}_{\alpha_\gamma}$ for α_γ . The three subgraphs occur only when \mathfrak{q} is of type $D_n(n-2)$. So, if \mathfrak{q} is not of type $D_n(n-2)$ then there are at most two subgraphs. In this case we denote by \mathfrak{l}_γ (resp. $\mathfrak{l}_{n\gamma}$) the simple ideal of \mathfrak{l} whose subgraph in the deleted Dynkin diagram contains (resp. does not contain) the node for α_γ . Thus the Levi subalgebra \mathfrak{l} may decompose into

$$(4.8) \quad \mathfrak{l} = \mathbb{C}H_q \oplus \mathfrak{l}_\gamma \oplus \mathfrak{l}_{n\gamma}.$$

Then, for the rest of this section, we assume that \mathfrak{q} is not of type $D_n(n-2)$, so that the Levi subalgebra \mathfrak{l} can be expressed as (4.8). Recall from Definition 4.3 that if \mathfrak{g} is of type \mathcal{T} then we say that the parabolic subalgebra \mathfrak{q} determined by $\alpha_i \in \Pi$ is of type $\mathcal{T}(i)$. Then the parabolic subalgebras \mathfrak{q} under consideration are given as follows:

$$(4.9) \quad B_n(i) \ (3 \leq i \leq n), \quad C_n(i) \ (2 \leq i \leq n-1), \quad D_n(i) \ (3 \leq i \leq n-3),$$

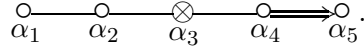
and

$$(4.10) \quad E_6(3), E_6(5), E_7(2), E_7(6), E_8(1), F_4(4).$$

Note that, in type A_n , any maximal parabolic subalgebra is of abelian type, and also that, in type G_2 , the two maximal parabolic subalgebras are either a 3-step or a Heisenberg parabolic subalgebra.

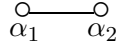
Write $\Pi(\mathfrak{l}_\gamma) = \{\alpha \in \Pi \mid \alpha \in \Delta(\mathfrak{l}_\gamma)\}$ and $\Pi(\mathfrak{l}_{n\gamma}) = \{\alpha \in \Pi \mid \alpha \in \Delta(\mathfrak{l}_{n\gamma})\}$. Example 4.11 below exhibits the subgraphs for \mathfrak{l}_γ and $\mathfrak{l}_{n\gamma}$ of \mathfrak{q} of type $B_5(3)$ with $\Pi(\mathfrak{l}_\gamma)$ and $\Pi(\mathfrak{l}_{n\gamma})$. One can find those data in Appendix A for each maximal parabolic subalgebra in (4.9) or (4.10).

Example 4.11. Let \mathfrak{q} be the parabolic subalgebra of type $B_5(3)$ with deleted Dynkin diagram

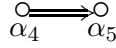


Observe that the unique simple root α_γ that is not orthogonal to the highest root γ is $\alpha_\gamma = \alpha_2$.

Therefore, the subgraph for \mathfrak{l}_γ is



and that for $\mathfrak{l}_{n\gamma}$ is



with $\Pi(\mathfrak{l}_\gamma) = \{\alpha_1, \alpha_2\}$ and $\Pi(\mathfrak{l}_{n\gamma}) = \{\alpha_4, \alpha_5\}$.

Remark 4.12. As α_γ is the unique simple root that is not orthogonal to γ , we have $\langle \gamma, \alpha_\gamma \rangle > 0$ and $\langle \gamma, \alpha \rangle = 0$ for any other simple roots α . In particular, $\langle \alpha, \gamma \rangle = 0$ for all $\alpha \in \Pi(\mathfrak{l}_{n\gamma})$.

4.4. The Highest Weights for \mathfrak{l}_γ , $\mathfrak{l}_{n\gamma}$, $\mathfrak{g}(1)$, and $\mathfrak{z}(\mathfrak{n})$. For the rest of this section we summarize technical lemmas on the L -highest weights for \mathfrak{l}_γ , $\mathfrak{l}_{n\gamma}$, $\mathfrak{g}(1)$, and $\mathfrak{z}(\mathfrak{n})$. These technical facts will be used in later computations.

Proposition 4.5 shows that $\mathfrak{z}(\mathfrak{n})$ has highest weight γ , which is the highest root of \mathfrak{g} . We denote by ξ_γ , $\xi_{n\gamma}$, and μ the highest weights for \mathfrak{l}_γ , $\mathfrak{l}_{n\gamma}$, and $\mathfrak{g}(1)$, respectively. These highest weights are summarized in Appendix A for each of the parabolic subalgebras under consideration. We remark that all these highest weights are indeed roots in Δ^+ . Observe that the highest weights ξ_γ and $\xi_{n\gamma}$ of \mathfrak{l}_γ and $\mathfrak{l}_{n\gamma}$, respectively, are also the highest roots of \mathfrak{l}_γ and $\mathfrak{l}_{n\gamma}$ as simple algebras; in particular, the multiplicities of $\alpha \in \Pi(\mathfrak{l}_\gamma)$ (resp. $\alpha \in \Pi(\mathfrak{l}_{n\gamma})$) in ξ_γ (resp. $\xi_{n\gamma}$) are all strictly positive.

Lemma 4.13. If α_q is the simple root that determines $\mathfrak{q} = \mathfrak{l} \oplus \mathfrak{g}(1) \oplus \mathfrak{z}(\mathfrak{n})$ then $\xi_\gamma + \alpha_q$ and $\xi_{n\gamma} + \alpha_q$ are roots.

Proof. We only prove that $\xi_\gamma + \alpha_q \in \Delta$; the other assertion that $\xi_{n\gamma} + \alpha_q \in \Delta$ can be proven similarly. It suffices to show that $\langle \xi_\gamma, \alpha_q \rangle < 0$, since both ξ_γ and α_q are roots. For $\alpha \in \Pi$ we observe that $\langle \alpha, \alpha_q \rangle < 0$ if α is adjacent to α_q in the Dynkin diagram and $\langle \alpha, \alpha_q \rangle = 0$ otherwise. An observation on the deleted Dynkin diagrams shows that there exists a unique simple root α_k in $\Pi(\mathfrak{l}_\gamma)$ that is adjacent to α_q . Since ξ_γ is the highest root for \mathfrak{l}_γ as a simple algebra, the multiplicity of α_k in ξ_γ is strictly positive. Thus $\langle \xi_\gamma, \alpha_q \rangle < 0$. \square

Lemma 4.14. If ξ_γ , $\xi_{n\gamma}$, μ , and γ are the highest weights of \mathfrak{l}_γ , $\mathfrak{l}_{n\gamma}$, $\mathfrak{g}(1)$, and $\mathfrak{z}(\mathfrak{n})$, respectively, then the following hold:

- (1) $\gamma - \xi_\gamma \in \Delta$, but $\gamma - \xi_{n\gamma} \notin \Delta$.
- (2) $\gamma - \mu \in \Delta$.
- (3) $\mu - \xi_\gamma, \mu - \xi_{n\gamma} \in \Delta$.

Proof. To prove $\gamma - \xi_{n\gamma} \notin \Delta$, observe that if n and m are the largest non-negative integers so that $\gamma - n\xi_{n\gamma} \in \Delta$ and $\gamma + m\xi_{n\gamma} \in \Delta$, respectively, then $\langle \gamma, \xi_{n\gamma}^\vee \rangle = n - m$. Since $\langle \gamma, \alpha^\vee \rangle = 0$ for all $\alpha \in \Delta(\mathfrak{l}_{n\gamma})$, we have $\langle \gamma, \xi_{n\gamma}^\vee \rangle = 0$ and so $n = m$. As $\xi_{n\gamma} \in \Delta^+$ and γ is the highest root, $\gamma + \xi_{n\gamma} \notin \Delta$. Therefore, $n = m = 0$, which concludes that $\gamma - \xi_{n\gamma}$ is not a root. To prove $\gamma - \xi_\gamma \in \Delta$, it suffices to show that $\langle \gamma, \xi_\gamma \rangle > 0$, since both γ and ξ_γ are roots. Write ξ_γ in terms of simple roots in $\Pi(\mathfrak{l}_\gamma)$. Observe that each $\alpha \in \Pi(\mathfrak{l}_\gamma)$ has positive multiplicity m_α in ξ_γ . As γ is orthogonal to α for any $\alpha \in \Pi(\mathfrak{l}_\gamma) \setminus \{\alpha_\gamma\}$, we have $\langle \gamma, \xi_\gamma \rangle = m_{\alpha_\gamma} \langle \gamma, \alpha_\gamma \rangle > 0$. The assertions (2) and (3) can be shown similarly. \square

The following technical lemma will simplify arguments concerning the long roots later. When \mathfrak{g} is simply laced, we regard any root as a long root.

Lemma 4.15. *Suppose that $\alpha \in \Delta$ is a long root. For any $\beta \in \Delta$, the following hold.*

- (1) *If $\beta - \alpha \in \Delta$ then $\langle \beta, \alpha^\vee \rangle = 1$.*
- (2) *If $\beta + \alpha \in \Delta$ then $\langle \beta, \alpha^\vee \rangle = -1$.*
- (3) *If $\beta \pm \alpha \in \Delta$ then $\beta \mp \alpha \notin \Delta$.*
- (4) *$\beta \pm 2\alpha \notin \Delta$.*

Proof. These simply follow from the standard arguments using the structure theory of Lie algebras. \square

Lemma 4.16. *If $\xi_\gamma, \xi_{n\gamma}, \mu$, and γ are the highest weights of $\mathfrak{l}_\gamma, \mathfrak{l}_{n\gamma}, \mathfrak{g}(1)$, and $\mathfrak{z}(\mathfrak{n})$, respectively, then the following hold:*

- (1) $\gamma - \mu + \xi_{n\gamma} \in \Delta$.
- (2) $\gamma - \mu - \xi_{n\gamma} \notin \Delta$.
- (3) *If ξ_γ is a long root then $\gamma - \mu \pm \xi_\gamma \notin \Delta$.*

Proof. Lemma 4.14 shows that $\gamma - \mu \in \Delta$. Then in order to prove (1), it is enough to show that $\langle \xi_{n\gamma}, \gamma - \mu \rangle < 0$. It follows from Remark 4.12 that $\langle \xi_{n\gamma}, \gamma \rangle = 0$. On the other hand, we have $\langle \xi_{n\gamma}, \mu \rangle > 0$ by the proof for (3) of Lemma 4.14. Therefore,

$$\langle \xi_{n\gamma}, \gamma - \mu \rangle = \langle \xi_{n\gamma}, \gamma \rangle - \langle \xi_{n\gamma}, \mu \rangle < 0.$$

When $\xi_{n\gamma}$ is a long root of \mathfrak{g} , the assertion (2) follows from (1) and Lemma 4.15. The data in Appendix A shows that $\xi_{n\gamma}$ is a long root unless \mathfrak{q} is of type $B_n(n-1)$. If \mathfrak{q} is of type $B_n(n-1)$ then we have $\gamma = \varepsilon_1 + \varepsilon_2$, $\mu = \varepsilon_1 + \varepsilon_n$, and $\xi_{n\gamma} = \varepsilon_n$. Thus $\gamma - \mu - \xi_{n\gamma} \notin \Delta$.

To show (3), observe that, by Lemma 4.14, we have $\gamma - \xi_\gamma, \mu - \xi_\gamma \in \Delta$. Since ξ_γ is assumed to be a long root, it follows from Lemma 4.15 that $\langle \gamma, \xi_\gamma^\vee \rangle = \langle \mu, \xi_\gamma^\vee \rangle = 1$. Therefore $\langle \gamma - \mu, \xi_\gamma^\vee \rangle = 0$, which forces that

$$(4.17) \quad \|\gamma - \mu \pm \xi_\gamma\|^2 = \|\gamma - \mu\|^2 + \|\xi_\gamma\|^2.$$

Since $\gamma - \mu$ is a root, we have $\|\gamma - \mu\| \neq 0$. As ξ_γ is assumed to be a long root, (4.17) implies that $(\gamma - \mu) \pm \xi_\gamma \notin \Delta$. \square

Remark 4.18. *Direct observation shows that ξ_γ is a long root, unless \mathfrak{q} is of type $C_n(i)$. If \mathfrak{q} is of type $C_n(i)$ then the data in Appendix A shows $\gamma = 2\varepsilon_1$, $\mu = \varepsilon_1 + \varepsilon_{i+1}$, and $\xi_\gamma = \varepsilon_1 - \varepsilon_i$. Thus $\gamma - \mu + \xi_\gamma \notin \Delta$, but $\gamma - \mu - \xi_\gamma \in \Delta$.*

5. THE Ω_1 SYSTEM

The aim of this section is to determine the complex parameter $s_1 \in \mathbb{C}$ for the line bundle \mathcal{L}_s so that the Ω_1 system of a maximal two-step nilpotent parabolic subalgebra \mathfrak{q} of non-Heisenberg type is conformally invariant on \mathcal{L}_{s_1} . To do so, it is essential to set up convenient normalizations.

If $\alpha, \beta \in \Delta$ then define

$$(5.1) \quad \begin{aligned} p_{\alpha,\beta} &= \max\{j \in \mathbb{Z}_{\geq 0} \mid \beta - j\alpha \in \Delta\} \text{ and} \\ q_{\alpha,\beta} &= \max\{j \in \mathbb{Z}_{\geq 0} \mid \beta + j\alpha \in \Delta\}. \end{aligned}$$

In particular, we have

$$(5.2) \quad \langle \beta, \alpha^\vee \rangle = p_{\alpha,\beta} - q_{\alpha,\beta}.$$

It is known that we can choose $X_\alpha \in \mathfrak{g}_\alpha$ and $H_\alpha \in \mathfrak{h}$ for each $\alpha \in \Delta$ in such a way that the following conditions hold (see for instance [6, Sections III.4 and III.5]). The reader may want to notice that our normalizations are different from those used in [1].

- (H1) For each $\alpha \in \Delta^+$, $\{X_\alpha, X_{-\alpha}, H_\alpha\}$ is an $\mathfrak{sl}(2, \mathbb{C})$ triple; in particular, we have $[X_\alpha, X_{-\alpha}] = H_\alpha$.
- (H2) For each $\alpha, \beta \in \Delta^+$, $[H_\alpha, X_\beta] = \beta(H_\alpha)X_\beta$.
- (H3) For $\alpha \in \Delta$ we have $\kappa(X_\alpha, X_{-\alpha}) = 1$.
- (H4) For $\alpha, \beta \in \Delta$ we have $\beta(H_\alpha) = \langle \alpha, \beta \rangle$.
- (H5) For $\alpha, \beta \in \Delta$ with $\alpha + \beta \neq 0$, there is a constant $N_{\alpha,\beta}$ so that

$$\begin{aligned} [X_\alpha, X_\beta] &= N_{\alpha,\beta} X_{\alpha+\beta} & \text{if } \alpha + \beta \in \Delta, \\ N_{\alpha,\beta} &= 0 & \text{if } \alpha + \beta \notin \Delta. \end{aligned}$$

- (H6) If $\alpha_1, \alpha_2, \alpha_3 \in \Delta^+$ with $\alpha_1 + \alpha_2 + \alpha_3 = 0$ then

$$N_{\alpha_1, \alpha_2} = N_{\alpha_2, \alpha_3} = N_{\alpha_3, \alpha_1}.$$

- (H7) If $\alpha, \beta \in \Delta$ and $\alpha + \beta \in \Delta$ then

$$N_{\alpha,\beta} N_{-\alpha,-\beta} = -\frac{q_{\alpha,\beta}(1 + p_{\alpha,\beta})}{2} \|\alpha\|^2$$

In particular, $N_{\alpha,\beta}$ is non-zero if $\alpha + \beta \in \Delta$.

We call the constants $N_{\alpha,\beta}$ structure constants.

As we have observed in Subsection 3.1, we use the covariant map τ_1 and the associated L -intertwining operators $\tilde{\tau}_1|_{V^*}$, where V^* are irreducible constituents of $\mathfrak{g}(-1)^* \otimes \mathfrak{g}(2)^*$. By Definition

3.1, the covariant map τ_1 is given by

$$\begin{aligned}\tau_1 : \mathfrak{g}(1) &\rightarrow \mathfrak{g}(-1) \otimes \mathfrak{z}(\mathfrak{n}) \\ X &\mapsto \text{ad}(X)\omega\end{aligned}$$

with $\omega = \sum_{\gamma_j \in \Delta(\mathfrak{z}(\mathfrak{n}))} X_{-\gamma_j} \otimes X_{\gamma_j}$. It is clear that τ_1 is not identically zero. Indeed, if $X = X_\mu$ with μ the highest weight for $\mathfrak{g}(1)$ then

$$\tau_1(X_\mu) = \text{ad}(X_\mu)\omega = \sum_{\Delta_\mu(\mathfrak{z}(\mathfrak{n}))} N_{\mu, -\gamma_j} X_{\mu-\gamma_j} \otimes X_{\gamma_j}$$

with $\Delta_\mu(\mathfrak{z}(\mathfrak{n})) = \{\gamma_j \in \Delta(\mathfrak{z}(\mathfrak{n})) \mid \mu - \gamma_j \in \Delta\}$. By Lemma 4.14, we have $\mu - \gamma \in \Delta$ with γ the highest weight for $\mathfrak{z}(\mathfrak{n})$, so $\Delta_\mu(\mathfrak{z}(\mathfrak{n})) \neq \emptyset$. Since the vectors $X_{\mu-\gamma_j} \otimes X_{\gamma_j}$ for $\gamma_j \in \Delta_\mu(\mathfrak{z}(\mathfrak{n}))$ are linearly independent, we have $\tau_1(X_\mu) \neq 0$.

For each irreducible constituent V^* of $\mathfrak{g}(-1)^* \otimes \mathfrak{z}(\mathfrak{n})^*$, there exists an associated L -intertwining operator $\tilde{\tau}_1|_{V^*} \in \text{Hom}_L(V^*, \mathcal{P}^1(\mathfrak{g}(1)))$ so that, for all $Y^* \in V^*$,

$$\tilde{\tau}_1|_{V^*}(Y^*)(X) = Y^*(\tau_1(X)).$$

Observe that the duality for V^* is defined with respect to the Killing form κ . Moreover, via the Killing form κ , we have $\mathfrak{g}(-1)^* \otimes \mathfrak{z}(\mathfrak{n})^* \cong \mathfrak{g}(1) \otimes \mathfrak{z}(\bar{\mathfrak{n}})$. Thus, if $Y^* = X_\alpha \otimes X_{-\gamma_t}$ with $\alpha \in \Delta(\mathfrak{g}(1))$ and $\gamma_t \in \Delta(\mathfrak{z}(\mathfrak{n}))$ then $Y^*(\tau_1(X))$ is given by

$$(5.3) \quad Y^*(\tau_1(X)) = \sum_{\gamma_j \in \Delta(\mathfrak{z}(\mathfrak{n}))} \kappa(X_\alpha, \text{ad}(X)X_{-\gamma_j}) \kappa(X_{-\gamma_t}, X_{\gamma_j}),$$

as $\tau_1(X) = \sum_{\gamma_j \in \Delta(\mathfrak{z}(\mathfrak{n}))} \text{ad}(X)X_{-\gamma_j} \otimes X_{\gamma_j}$.

Now we wish to determine all the irreducible constituents V^* of $\mathfrak{g}(1) \otimes \mathfrak{z}(\bar{\mathfrak{n}})$, so that $\tilde{\tau}_1|_{V^*}$ are not identically zero. Observe that $\mathcal{P}^1(\mathfrak{g}(1)) \cong \text{Sym}^1(\mathfrak{g}(-1)) = \mathfrak{g}(-1)$ and that $\mathfrak{g}(-1)$ is an irreducible L -module, as \mathfrak{q} is a maximal parabolic subalgebra. Thus, if $\tilde{\tau}_1|_{V^*}$ is not identically zero then $V^* \cong \mathfrak{g}(-1)$. Proposition 5.4 below shows that the converse also holds.

Proposition 5.4. *Let V^* be an irreducible constituent of $\mathfrak{g}(1) \otimes \mathfrak{z}(\bar{\mathfrak{n}})$. Then $\tilde{\tau}_1|_{V^*}$ is not identically zero if and only if $V^* \cong \mathfrak{g}(-1)$.*

Proof. First observe that $\mathfrak{g}(-1)$ is an irreducible constituent of $\mathfrak{g}(1) \otimes \mathfrak{z}(\bar{\mathfrak{n}})$. Indeed, since τ_1 is linear, we have $\tau_1(\mathfrak{g}(1)) \cong \mathfrak{g}(1)$ as an L -module; in particular, $\mathfrak{g}(1)$ is an irreducible constituent of $\mathfrak{g}(-1) \otimes \mathfrak{z}(\mathfrak{n})$. Therefore $\mathfrak{g}(-1) \cong \mathfrak{g}(1)^*$ is an irreducible constituent of $\mathfrak{g}(1) \otimes \mathfrak{z}(\bar{\mathfrak{n}}) \cong (\mathfrak{g}(-1) \otimes \mathfrak{z}(\mathfrak{n}))^*$.

To prove $\tilde{\tau}_1|_{\mathfrak{g}(-1)}$ is a non-zero map, it suffices to show that $\tilde{\tau}_1|_{\mathfrak{g}(-1)}(Y^*) \neq 0$ for some $Y^* \in \mathfrak{g}(-1) \subset \mathfrak{g}(1) \otimes \mathfrak{z}(\bar{\mathfrak{n}})$. To do so, consider a map

$$\begin{aligned}\bar{\tau}_1 : \mathfrak{g}(-1) &\rightarrow \mathfrak{g}(1) \otimes \mathfrak{z}(\bar{\mathfrak{n}}) \\ \bar{X} &\mapsto \text{ad}(\bar{X})\bar{\omega}\end{aligned}$$

with $\bar{\omega} = \sum_{\gamma_t \in \Delta(\mathfrak{z}(\mathfrak{n}))} X_{\gamma_t} \otimes X_{-\gamma_t}$. This is a non-zero L -intertwining operator. Thus $\bar{\tau}_1(\mathfrak{g}(-1)) \cong \mathfrak{g}(-1)$ as an L -module, and $\bar{\tau}_1(X_{-\alpha})$ is a weight vector with weight $-\alpha$ for all $\alpha \in \Delta(\mathfrak{g}(1))$. As $\mathfrak{g}(1)$ has highest weight μ , the lowest weight for $\mathfrak{g}(-1)$ is $-\mu$.

Now we set

$$c_\mu = \sum_{\gamma_t \in \Delta_\mu(\mathfrak{z}(\mathfrak{n}))} N_{-\mu, \gamma_t} N_{\mu, -\gamma_t}$$

with $\Delta_\mu(\mathfrak{z}(\mathfrak{n})) = \{\gamma_t \in \Delta(\mathfrak{z}(\mathfrak{n})) \mid \gamma_t - \mu \in \Delta\}$. By Lemma 4.14, it follows that $\gamma - \mu \in \Delta$; in particular, $\Delta_\mu(\mathfrak{z}(\mathfrak{n})) \neq \emptyset$. The normalization (H7) shows that $N_{-\mu, \gamma_t} N_{\mu, -\gamma_t} < 0$ for all $\gamma_t \in \Delta_\mu(\mathfrak{z}(\mathfrak{n}))$. Therefore $c_\mu \neq 0$. Then define $Y_l^* \in \mathfrak{g}(-1)$ by means of

$$Y_l^* = \frac{1}{c_\mu} \bar{\tau}_1(X_{-\mu}) = \frac{1}{c_\mu} \sum_{\gamma_t \in \Delta_\mu(\mathfrak{z}(\mathfrak{n}))} N_{-\mu, \gamma_t} X_{\gamma_t - \mu} \otimes X_{-\gamma_t}.$$

We claim that $\tilde{\tau}_1|_{\mathfrak{g}(-1)}(Y_l^*)(X) \neq 0$. By (5.3), the polynomial $\tilde{\tau}_1|_{\mathfrak{g}(-1)}(Y_l^*)(X)$ is

$$\begin{aligned} \tilde{\tau}_1|_{\mathfrak{g}(-1)}(Y_l^*)(X) &= Y_l^*(\tau_1(X)) \\ &= \frac{1}{c_\mu} \sum_{\substack{\gamma_t \in \Delta_\mu(\mathfrak{z}(\mathfrak{n})) \\ \gamma_j \in \Delta(\mathfrak{z}(\mathfrak{n}))}} N_{-\mu, \gamma_t} \kappa(X_{\gamma_t - \mu}, \text{ad}(X)X_{-\gamma_j}) \kappa(X_{-\gamma_t}, X_{\gamma_j}) \\ &= \frac{1}{c_\mu} \sum_{\gamma_t \in \Delta_\mu(\mathfrak{z}(\mathfrak{n}))} N_{-\mu, \gamma_t} \kappa(X_{\gamma_t - \mu}, \text{ad}(X)X_{-\gamma_t}). \end{aligned}$$

Write $X = \sum_{\alpha \in \Delta(\mathfrak{g}(1))} \eta_\alpha X_\alpha$, where $\eta_\alpha \in \mathfrak{n}^*$ is the coordinate dual to X_α with respect to the Killing form κ . Then,

$$\begin{aligned} \tilde{\tau}_1|_{\mathfrak{g}(-1)}(Y_l^*)(X) &= \frac{1}{c_\mu} \sum_{\gamma_t \in \Delta_\mu(\mathfrak{z}(\mathfrak{n}))} N_{-\mu, \gamma_t} \kappa(X_{\gamma_t - \mu}, \text{ad}(X)X_{-\gamma_t}) \\ &= \frac{1}{c_\mu} \sum_{\substack{\alpha \in \Delta(\mathfrak{g}(1)) \\ \gamma_t \in \Delta_\mu(\mathfrak{z}(\mathfrak{n}))}} N_{-\mu, \gamma_t} \eta_\alpha \kappa(X_{\gamma_t - \mu}, \text{ad}(X_\alpha)X_{-\gamma_t}) \\ &= \frac{1}{c_\mu} \sum_{\gamma_t \in \Delta_\mu(\mathfrak{z}(\mathfrak{n}))} N_{-\mu, \gamma_t} N_{\mu, -\gamma_t} \eta_\mu \\ &= \eta_\mu \\ (5.5) \quad &= \kappa(X, X_{-\mu}). \end{aligned}$$

Hence $\tilde{\tau}_1|_{\mathfrak{g}(-1)}(Y_l^*)(X) \neq 0$. □

Since only $\mathfrak{g}(-1)$ contributes to the construction of the Ω_1 systems, we simply refer to the Ω_1 system as the $\Omega_1|_{\mathfrak{g}(-1)}$ system. As we observed in Subsection 3.1, the operator $\Omega_1|_{\mathfrak{g}(-1)} : \mathfrak{g}(-1) \rightarrow \mathbb{D}(\mathcal{L}_s)^{\bar{\mathfrak{n}}}$ is obtained via the composition of maps

$$\mathfrak{g}(-1) \xrightarrow{\tilde{\tau}_1|_{\mathfrak{g}(-1)}} \mathcal{P}^1(\mathfrak{g}(1)) \rightarrow \mathfrak{g}(-1) \xrightarrow{\sigma} \mathcal{U}(\bar{\mathfrak{n}}) \xrightarrow{R} \mathbb{D}(\mathcal{L}_s)^{\bar{\mathfrak{n}}}.$$

By (5.5), we have $\tilde{\tau}_1|_{\mathfrak{g}(-1)}(Y_l^*)(X) = \kappa(X, X_{-\mu})$. Therefore,

$$\Omega_1(Y_l^*) = R(X_{-\mu}).$$

Now, for all $\alpha \in \Delta(\mathfrak{g}(1))$, set

$$Y_{-\alpha} = \bar{\tau}_1(X_{-\alpha}).$$

Then, as $Y_l^* = (1/c_\mu)\bar{\tau}_1(X_{-\mu})$, we have $\Omega_1(Y_{-\mu}) = c_\mu R(X_{-\mu})$. Since both $\Omega_1|_{\mathfrak{g}(-1)}$ and $\bar{\tau}_1$ are L_0 -intertwining operators and since $\mathfrak{g}(-1) = \mathcal{U}(\mathfrak{l})X_{-\mu}$, for any $\alpha \in \Delta(\mathfrak{g}(1))$, we obtain

$$(5.6) \quad \Omega_1(Y_{-\alpha}) = c_\alpha R(X_{-\alpha})$$

with some constant c_α . Thus, if $\Delta(\mathfrak{g}(1)) = \{\alpha_1, \dots, \alpha_m\}$ then the Ω_1 system is given by

$$R(X_{-\alpha_1}), \dots, R(X_{-\alpha_m}).$$

Theorem 5.7. *Let \mathfrak{g} be a complex simple Lie algebra, and let \mathfrak{q} be a maximal two-step nilpotent parabolic subalgebra of non-Heisenberg type. Then the Ω_1 system is conformally invariant on \mathcal{L}_s if and only if $s = 0$.*

Proof. By Remark 3.8, we only need to show that the condition (S2) in Definition 2.3 holds if and only if $s = 0$. By Theorem 2.17, for any $Y \in \mathfrak{g}$ and any $f \in C^\infty(\bar{N}_0, \mathbb{C}_{\chi^s})$, we have

$$\begin{aligned} & ([\pi_s(Y), R(X_{-\alpha_j})] \bullet f)(\bar{n}) \\ &= (R([(Ad(\bar{n}^{-1})Y)_{\mathfrak{q}}, X_{-\alpha_j}]_{\bar{n}}] \bullet f)(\bar{n}) - s\lambda_{\mathfrak{q}}([Ad(\bar{n}^{-1})Y, X_{-\alpha_j}]_{\mathfrak{q}})f(\bar{n}). \end{aligned}$$

Hence, the condition (S2) holds if and only if $s = 0$. \square

6. SPECIAL CONSTITUENTS OF $\mathfrak{l} \otimes \mathfrak{z}(\mathfrak{n})$

Our next goal is to construct the Ω_2 systems and to find their special values. To do so, we need to detect the irreducible constituents V^* of $\mathfrak{l}^* \otimes \mathfrak{z}(\mathfrak{n})^*$ so that $\tilde{\tau}_2|_{V^*}$ is not identically zero. (See Subsection 3.1 for the general construction of the Ω_k systems.) In this section we shall show preliminary results to find such irreducible constituents.

6.1. Irreducible Decomposition of $\mathfrak{l}_\gamma \otimes \mathfrak{z}(\mathfrak{n})$. We continue with $\mathfrak{q} = \mathfrak{l} \oplus \mathfrak{g}(1) \oplus \mathfrak{z}(\mathfrak{n})$ a maximal two-step nilpotent parabolic subalgebra of non-Heisenberg type listed in (4.9) or (4.10), and $Q = LN = N_G(\mathfrak{q})$. The Levi subgroup L acts on $\mathfrak{l} \otimes \mathfrak{z}(\mathfrak{n}) \subset \mathfrak{g} \otimes \mathfrak{g}$ via the standard action on the tensor product induced by the adjoint representation on \mathfrak{l} and $\mathfrak{z}(\mathfrak{n})$. As L is complex reductive, this action is completely reducible. Since $\mathfrak{l} = \mathfrak{z}(\mathfrak{l}) \oplus \mathfrak{l}_\gamma \oplus \mathfrak{l}_{n\gamma}$ with $\mathfrak{z}(\mathfrak{l}) = \mathbb{C}H_{\mathfrak{q}}$, we have

$$(6.1) \quad \mathfrak{l} \otimes \mathfrak{z}(\mathfrak{n}) = (\mathbb{C}H_{\mathfrak{q}} \otimes \mathfrak{z}(\mathfrak{n})) \oplus (\mathfrak{l}_\gamma \otimes \mathfrak{z}(\mathfrak{n})) \oplus (\mathfrak{l}_{n\gamma} \otimes \mathfrak{z}(\mathfrak{n})).$$

It is clear that $\mathbb{C}H_{\mathfrak{q}} \otimes \mathfrak{z}(\mathfrak{n}) \cong \mathfrak{z}(\mathfrak{n}) = \mathfrak{g}(2)$ as an L -module. Thus, by Corollary 4.5, $\mathbb{C}H_{\mathfrak{q}} \otimes \mathfrak{z}(\mathfrak{n})$ is L -irreducible. It is also easy to show that $\mathfrak{l}_{n\gamma} \otimes \mathfrak{z}(\mathfrak{n})$ is L -irreducible. Let L_γ (resp. $L_{n\gamma}$) be the analytic subgroup of L with Lie algebra \mathfrak{l}_γ (resp. $\mathfrak{l}_{n\gamma}$). As in Subsection 4.2, we call a weight ν for a finite dimensional L -module V a highest weight for V if the restriction $\nu|_{\mathfrak{h}_{ss}}$ onto \mathfrak{h}_{ss} is a highest weight for V as an L_{ss} -module.

Proposition 6.2. *Suppose that $\mathfrak{l}_{n\gamma} \neq 0$. If $\xi_{n\gamma}$ and γ are the highest weights of $\mathfrak{l}_{n\gamma}$ and $\mathfrak{z}(\mathfrak{n})$, respectively, then $\mathfrak{l}_{n\gamma} \otimes \mathfrak{z}(\mathfrak{n})$ is the irreducible L -module with highest weight $\xi_{n\gamma} + \gamma$.*

Proof. First we observe that $L_{n\gamma}$ acts trivially on $\mathfrak{z}(\mathfrak{n})$. By Corollary 4.5, we have $\mathfrak{z}(\mathfrak{n}) = \mathfrak{g}(2) = \mathcal{U}([\mathfrak{l}, \mathfrak{l}])X_\gamma$. By the observation made in Remark 4.12, it follows that $\alpha \perp \gamma$ for all $\alpha \in \Delta(\mathfrak{l}_{n\gamma})$. Thus $\mathfrak{z}(\mathfrak{n}) = \mathcal{U}(\mathfrak{l}_\gamma)X_\gamma$. Hence $L_{n\gamma}$ acts trivially; in particular, the irreducible L -module $\mathfrak{z}(\mathfrak{n})$ is L_γ -irreducible. On the other hand, it is clear that L_γ acts on $\mathfrak{l}_{n\gamma}$ trivially. Therefore the representation

$(L, \text{Ad} \otimes \text{Ad}, \mathfrak{l}_{n\gamma} \otimes \mathfrak{z}(\mathfrak{n}))$ is equivalent to $(L_\gamma \times L_{n\gamma}, \text{Ad} \hat{\otimes} \text{Ad}, \mathfrak{l}_{n\gamma} \otimes \mathfrak{z}(\mathfrak{n}))$, where $\hat{\otimes}$ denotes the outer tensor product. Since $\mathfrak{l}_{n\gamma}$ and $\mathfrak{z}(\mathfrak{n})$ have highest weight $\xi_{n\gamma}$ and γ , respectively, the lemma follows. \square

Now we focus on the irreducible decomposition of $\mathfrak{l}_\gamma \otimes \mathfrak{z}(\mathfrak{n})$. As noted in the proof for Proposition 6.2, the subgroup $L_{n\gamma}$ acts trivially on $\mathfrak{l}_\gamma \otimes \mathfrak{z}(\mathfrak{n})$. Hence we study $\mathfrak{l}_\gamma \otimes \mathfrak{z}(\mathfrak{n})$ as an L_γ -module. For $\lambda \in \mathfrak{h}^*$ with $\langle \lambda, \alpha^\vee \rangle \in \mathbb{Z}_{\geq 0}$ for all $\alpha \in \Pi(\mathfrak{l}_\gamma)$, we will denote by $V(\lambda)$ the irreducible constituent with highest weight $\lambda|_{\mathfrak{h}_\gamma}$, where $\mathfrak{h}_\gamma = \mathfrak{h} \cap \mathfrak{l}_\gamma$. For classical algebras, we use the standard realization of the roots ε_i , the dual basis of the standard orthonormal basis for \mathbb{R}^n . For exceptional algebras the Bourbaki conventions are used to label the simple roots.

Theorem 6.3. *The L -module $\mathfrak{l}_\gamma \otimes \mathfrak{z}(\mathfrak{n})$ is reducible. If $V(\lambda)$ denotes the irreducible representation of L with highest weight $\lambda|_{\mathfrak{h}_\gamma}$ then the irreducible decomposition of $\mathfrak{l}_\gamma \otimes \mathfrak{z}(\mathfrak{n})$ is given as follows.*

(1) $B_n(i)$, $3 \leq i \leq n$:

$$\begin{cases} V(\xi_\gamma + \gamma) \oplus V(\gamma) \oplus V(\xi_\gamma + (\varepsilon_1 + \varepsilon_3)) & \text{if } i = 3 \\ V(\xi_\gamma + \gamma) \oplus V(\gamma) \oplus V(\xi_\gamma + (\varepsilon_1 + \varepsilon_i)) \oplus V(\xi_\gamma + (\varepsilon_2 + \varepsilon_3)) & \text{if } 4 \leq i \leq n \end{cases}$$

(2) $C_n(i)$, $2 \leq i \leq n-1$:

$$\begin{cases} V(\xi_\gamma + \gamma) \oplus V(\gamma) \oplus V(\xi_\gamma + 2\varepsilon_2) & \text{if } i = 2 \\ V(\xi_\gamma + \gamma) \oplus V(\gamma) \oplus V(\xi_\gamma + (\varepsilon_2 + \varepsilon_i)) \oplus V(\xi_\gamma + (\varepsilon_1 + \varepsilon_2)) & \text{if } 3 \leq i \leq n-1 \end{cases}$$

(3) $D_n(i)$, $3 \leq i \leq n-3$:

$$\begin{cases} V(\xi_\gamma + \gamma) \oplus V(\gamma) \oplus V(\xi_\gamma + (\varepsilon_1 + \varepsilon_3)) & \text{if } i = 3 \\ V(\xi_\gamma + \gamma) \oplus V(\gamma) \oplus V(\xi_\gamma + (\varepsilon_1 + \varepsilon_i)) \oplus V(\xi_\gamma + (\varepsilon_2 + \varepsilon_3)) & \text{if } 4 \leq i \leq n-3 \end{cases}$$

(4) All exceptional cases ($E_6(3)$, $E_6(5)$, $E_7(2)$, $E_7(6)$, $E_8(1)$, $F_4(4)$):

$$V(\xi_\gamma + \gamma) \oplus V(\gamma) \oplus V(\xi_\gamma + \gamma_0),$$

where γ_0 is the following root contributing to $\mathfrak{z}(\mathfrak{n})$:

$$E_6(3) : \gamma_0 = \alpha_1 + \alpha_2 + 2\alpha_3 + 3\alpha_4 + 2\alpha_5 + \alpha_6$$

$$E_6(5) : \gamma_0 = \alpha_1 + \alpha_2 + 2\alpha_3 + 3\alpha_4 + 2\alpha_5 + \alpha_6$$

$$E_7(2) : \gamma_0 = \alpha_1 + 2\alpha_2 + 3\alpha_3 + 4\alpha_4 + 3\alpha_5 + 2\alpha_6 + \alpha_7$$

$$E_7(6) : \gamma_0 = \alpha_1 + 2\alpha_2 + 2\alpha_3 + 4\alpha_4 + 3\alpha_5 + 2\alpha_6 + \alpha_7$$

$$E_8(1) : \gamma_0 = 2\alpha_1 + 3\alpha_2 + 4\alpha_3 + 6\alpha_4 + 5\alpha_5 + 4\alpha_6 + 2\alpha_7 + \alpha_8$$

$$F_4(4) : \gamma_0 = \alpha_1 + 2\alpha_2 + 4\alpha_3 + 2\alpha_4.$$

Proof. To prove this theorem we just use the standard character formula due to Klimyk ([11, Corollary]) for $\mathfrak{l}_\gamma \otimes \mathfrak{z}(\mathfrak{n})$. For the details, see Chapter 5 of [15]. \square

6.2. Special Constituents. Given irreducible constituent $V(\nu)$ in $\mathfrak{l} \otimes \mathfrak{z}(\mathfrak{n})$, we build L -intertwining map

$$\tilde{\tau}_2|_{V(\nu)^*} \in \text{Hom}_L(V(\nu)^*, \mathcal{P}^2(\mathfrak{g}(1)))$$

with $V(\nu)^*$ the dual of $V(\nu)$ with respect to the Killing form κ . From $\tilde{\tau}_2|_{V(\nu)^*}$, we construct operator $\Omega_2|_{V(\nu)^*} : V(\nu)^* \rightarrow \mathbb{D}(\mathcal{L}_s)^{\bar{n}}$. To do so, it is necessary to determine which irreducible constituents

$V(\nu)$ have the property that $\tilde{\tau}_2|_{V(\nu)^*} \neq 0$. Thus, next, by using the above decomposition results, we shall determine such irreducible constituents.

First we observe the vector space isomorphism $\mathcal{P}^2(\mathfrak{g}(1)) \cong \text{Sym}^2(\mathfrak{g}(1))^*$. With the natural L -action on $\mathcal{P}^2(\mathfrak{g}(1))$ and $\text{Sym}^2(\mathfrak{g}(1))^*$, this vector space isomorphism is L -equivariant. Thus, if $\tilde{\tau}_2|_{V(\nu)^*}$ is a non-zero map then $V(\nu)$ is an irreducible constituent of $\text{Sym}^2(\mathfrak{g}(1)) \subset \mathfrak{g}(1) \otimes \mathfrak{g}(1)$; in particular, the weight ν is of the form $\nu = \mu + \epsilon$ for some $\epsilon \in \Delta(\mathfrak{g}(1))$, where μ is the highest weight of $\mathfrak{g}(1)$.

One can see from the decompositions in Theorem 6.3 that $V(\gamma)$ is an irreducible constituent of $\mathfrak{l} \otimes \mathfrak{z}(\mathfrak{n})$ for any \mathfrak{q} under consideration. By Lemma 4.14, we have $\gamma = \mu + \epsilon$ for some $\epsilon \in \Delta(\mathfrak{g}(1))$. Now we claim that $\tilde{\tau}_2|_{V(\gamma)^*}$ is identically zero. It is well-known that

$$(6.4) \quad \mathfrak{g}(1) \otimes \mathfrak{g}(1) = \text{Sym}^2(\mathfrak{g}(1)) \oplus \wedge^2(\mathfrak{g}(1))$$

as an L -module. Since each weight space for $\mathfrak{g}(1)$ is one-dimensional as weights for $\mathfrak{g}(1)$ are roots of \mathfrak{g} , the L -module decomposition (6.4) is multiplicity free.

Proposition 6.5. *The L -module $V(\gamma)$ is an irreducible constituent of $\wedge^2(\mathfrak{g}(1))$.*

Proof. Define a linear map $\varphi : \mathfrak{z}(\mathfrak{n}) \rightarrow \wedge^2(\mathfrak{g}(1))$ by means of

$$\varphi(W) = \sum_{\beta \in \Delta(\mathfrak{g}(1))} \text{ad}(W)X_{-\beta} \wedge X_{\beta}.$$

By using an argument similar to that for Lemma 3.3, one can show that φ is L -equivariant. Then, since $\mathfrak{z}(\mathfrak{n}) \cong V(\gamma)$ as an irreducible L -module, it suffices to show that φ is a non-zero map. Write $\Delta_{\gamma}(\mathfrak{g}(1)) = \{\beta \in \Delta(\mathfrak{g}(1)) \mid \gamma - \beta \in \Delta\}$. By Lemma 4.14, we have $\gamma - \mu \in \Delta$. Hence $\Delta_{\gamma}(\mathfrak{g}(1)) \neq \emptyset$. By writing $\beta' = \gamma - \beta$ for $\beta \in \Delta_{\gamma}(\mathfrak{g}(1))$, $\varphi(X_{\gamma})$ is given by

$$\varphi(X_{\gamma}) = \sum_{\beta \in \Delta(\mathfrak{g}(1))} \text{ad}(X_{\gamma})X_{-\beta} \wedge X_{\beta} = \sum_{\beta \in \Delta_{\gamma}(\mathfrak{g}(1))} N_{\gamma, -\beta} X_{\beta'} \wedge X_{\beta}.$$

Observe that for each $\beta \in \Delta_{\gamma}(\mathfrak{g}(1))$, we have $\gamma - \beta \in \Delta_{\gamma}(\mathfrak{g}(1))$. Moreover, by the normalization (H6) of our normalizations in Section 5, it follows that $N_{\gamma, -\beta'} = -N_{\gamma, -\beta}$. Therefore,

$$(6.6) \quad N_{\gamma, -\beta} X_{\beta'} \wedge X_{\beta} + N_{\gamma, -\beta'} X_{\beta} \wedge X_{\beta'} = 2N_{\gamma, -\beta} X_{\beta'} \wedge X_{\beta}.$$

Since $N_{\gamma, -\beta} \neq 0$ for $\beta \in \Delta_{\gamma}(\mathfrak{g}(1))$, equation (6.6) is non-zero. On the other hand, if $\beta \in \Delta_{\gamma}(\mathfrak{g}(1))$ and $\eta \in \Delta_{\gamma}(\mathfrak{g}(1))$ is so that $\eta \neq \beta, \beta'$ then $X_{\beta'} \wedge X_{\beta}$ and $X_{\eta} \wedge X_{\beta}$ are linearly independent. Hence, $\varphi(X_{\gamma}) \neq 0$. \square

Definition 6.7. *An irreducible constituent $V(\nu)$ of $\mathfrak{l} \otimes \mathfrak{z}(\mathfrak{n})$ is called **special** if $\nu \neq \gamma$ and there exists $\epsilon \in \Delta(\mathfrak{g}(1))$ so that $\nu = \mu + \epsilon$, where μ and γ are the highest weights for $\mathfrak{g}(1)$ and $\mathfrak{z}(\mathfrak{n})$, respectively.*

Proposition 6.8. *Let $V(\nu)$ be an irreducible constituent of $\mathfrak{l} \otimes \mathfrak{z}(\mathfrak{n})$. Then $\tilde{\tau}_2|_{V(\nu)^*}$ is not identically zero only if $V(\nu)$ is a special constituent of $\mathfrak{l} \otimes \mathfrak{z}(\mathfrak{n})$.*

Proof. At the beginning of this section we observed that if $\tilde{\tau}_2|_{V(\nu)^*} \neq 0$ then ν must be of the form $\nu = \mu + \epsilon$ for some $\epsilon \in \Delta(\mathfrak{g}(1))$. Then $V(\nu)$ is either a special constituent or $V(\gamma)$ (by Lemma

4.14, γ satisfies the form). However, by Proposition 6.5, it follows that $\tilde{\tau}_2|_{V(\gamma)^*}$ is identically zero. Therefore, $V(\nu)$ must be a special constituent. \square

Now we determine all the special constituents of $\mathfrak{l} \otimes \mathfrak{z}(\mathfrak{n})$. Since $\mathfrak{l} \otimes \mathfrak{z}(\mathfrak{n}) = (\mathbb{C}H_{\mathfrak{q}} \otimes \mathfrak{z}(\mathfrak{n})) \oplus ([\mathfrak{l}, \mathfrak{l}] \otimes \mathfrak{z}(\mathfrak{n}))$ and $\mathbb{C}H_{\mathfrak{q}} \otimes \mathfrak{z}(\mathfrak{n}) = V(\gamma)$, it suffices to consider $[\mathfrak{l}, \mathfrak{l}] \otimes \mathfrak{z}(\mathfrak{n}) = (\mathfrak{l}_{\gamma} \otimes \mathfrak{z}(\mathfrak{n})) \oplus (\mathfrak{l}_{n\gamma} \otimes \mathfrak{z}(\mathfrak{n}))$. We start by observing that, by Proposition 6.2, $\mathfrak{l}_{n\gamma} \otimes \mathfrak{z}(\mathfrak{n}) = V(\xi_{n\gamma} + \gamma)$.

Proposition 6.9. *Suppose that $\mathfrak{l}_{n\gamma} \neq 0$. Then the irreducible constituent $V(\xi_{n\gamma} + \gamma)$ is special.*

Proof. We need to show that $\xi_{n\gamma} + \gamma = \mu + \beta$ for some $\beta \in \Delta(\mathfrak{g}(1))$. This is precisely the statement (1) of Lemma 4.16. \square

We next consider the constituent $V(\xi_{\gamma} + \gamma)$ of $\mathfrak{l}_{\gamma} \otimes \mathfrak{z}(\mathfrak{n}) = V(\xi_{\gamma}) \otimes V(\gamma)$.

Lemma 6.10. *The irreducible constituent $V(\xi_{\gamma} + \gamma)$ of $\mathfrak{l}_{\gamma} \otimes \mathfrak{z}(\mathfrak{n})$ is not special.*

Proof. Lemma 4.16 and Remark 4.18 show that $\xi_{\gamma} + \gamma - \mu \notin \Delta(\mathfrak{g}(1))$, which implies that $\xi_{\gamma} + \gamma \neq \mu + \beta$ for all $\beta \in \Delta(\mathfrak{g}(1))$. \square

We determine all the special constituents of $\mathfrak{l}_{\gamma} \otimes \mathfrak{z}(\mathfrak{n})$ in two steps. First we assume that \mathfrak{g} is a classical algebra, and then consider the case that \mathfrak{g} is an exceptional algebra.

For classical cases the parabolic subalgebras \mathfrak{q} under consideration are of type $B_n(i)$ ($3 \leq i \leq n$), $C_n(i)$ ($2 \leq i \leq n-1$), or $D_n(i)$ ($3 \leq i \leq n-3$). It will be convenient to write $\beta \in \Delta(\mathfrak{g}(1))$ in terms of the fundamental weights of \mathfrak{l}_{γ} and $\mathfrak{l}_{n\gamma}$. It is clear from the deleted Dynkin diagrams that, for each of the cases, $\Pi(\mathfrak{l}_{\gamma})$ and $\Pi(\mathfrak{l}_{n\gamma})$ are given by

$$\Pi(\mathfrak{l}_{\gamma}) = \{\alpha_r \mid 1 \leq r \leq i-1\} \quad \text{and} \quad \Pi(\mathfrak{l}_{n\gamma}) = \{\alpha_{i+s} \mid 1 \leq s \leq n-i\},$$

where α_j are the simple roots with the standard numbering. By using the standard realizations of roots, we have $\alpha_r = \varepsilon_r - \varepsilon_{r+1}$ for $1 \leq r \leq i-1$, $\alpha_{i+s} = \varepsilon_{i+s} - \varepsilon_{i+s+1}$ for $1 \leq s \leq n-i-1$, and

$$\alpha_n = \begin{cases} \varepsilon_n & \text{if } \mathfrak{g} \text{ is of type } B_n \\ 2\varepsilon_n & \text{if } \mathfrak{g} \text{ is of type } C_n \\ \varepsilon_{n-1} + \varepsilon_n & \text{if } \mathfrak{g} \text{ is of type } D_n. \end{cases}$$

The data in Appendix A shows that if \mathfrak{q} is of type $B_n(i)$ then

$$\Delta(\mathfrak{g}(1)) = \{\varepsilon_j \pm \varepsilon_k \mid 1 \leq j \leq i \text{ and } i+1 \leq k \leq n\} \cup \{\varepsilon_j \mid 1 \leq j \leq i\}$$

and if \mathfrak{q} is of type $C_n(i)$ or $D_n(i)$ then

$$\Delta(\mathfrak{g}(1)) = \{\varepsilon_j \pm \varepsilon_k \mid 1 \leq j \leq i \text{ and } i+1 \leq k \leq n\}.$$

Since we have two simple algebras \mathfrak{l}_{γ} and $\mathfrak{l}_{n\gamma}$, we use the notation ϖ_r for the fundamental weights of $\alpha_r \in \Pi(\mathfrak{l}_{\gamma})$ and $\tilde{\varpi}_s$ for those of $\alpha_{i+s} \in \Pi(\mathfrak{l}_{n\gamma})$. Direct computation then shows that each $\beta \in \Delta(\mathfrak{g}(1))$ is exactly one of the following forms:

$$(6.11) \quad \beta = \begin{cases} \varpi_1 + \sum_{s=1}^{n-i} \tilde{m}_s \tilde{\varpi}_s, \\ (-\varpi_r + \varpi_{r+1}) + \sum_{s=1}^{n-i} \tilde{m}_s \tilde{\varpi}_s \text{ with } 1 \leq r \leq i-2, \text{ or} \\ -\varpi_{i-1} + \sum_{s=1}^{n-i} \tilde{m}_s \tilde{\varpi}_s \end{cases}$$

for some $\tilde{m}_s \in \mathbb{Z}$.

Proposition 6.12. *Let $V(\nu)$ be an irreducible constituent of $\mathfrak{l}_\gamma \otimes \mathfrak{z}(\mathfrak{n})$.*

- (1) *If \mathfrak{q} is of type $B_n(i)$ ($3 \leq i \leq n$) or $D_n(i)$ ($3 \leq i \leq n-3$) then $V(\nu)$ is a special constituent if and only if $\nu = 2\varepsilon_1$.*
- (2) *If \mathfrak{q} is of type $C_n(i)$ ($2 \leq i \leq n-1$) then $V(\nu)$ is a special constituent if and only if $\nu = \varepsilon_1 + \varepsilon_2$.*

Proof. Suppose that \mathfrak{q} is of type $B_n(i)$, $C_n(i)$, or $D_n(i)$. By Definition 6.7, we need to find all ν of the form $\nu = \mu + \beta$ for some $\beta \in \Delta(\mathfrak{g}(1))$. Here μ , the highest weight for $\mathfrak{g}(1)$, is

$$\mu = \begin{cases} \varepsilon_1 + \varepsilon_{i+1} & \text{if } \mathfrak{q} \text{ is of type } B_n(i) \text{ with } i \neq n, C_n(i), \text{ or } D_n(i) \\ \varepsilon_1 & \text{if } \mathfrak{q} \text{ is of type } B_n(n). \end{cases}$$

We write μ in terms of the fundamental weights of \mathfrak{l}_γ and $\mathfrak{l}_{n\gamma}$; that is,

$$(6.13) \quad \mu = \begin{cases} \varpi_1 + \tilde{\varpi}_1 & \text{if } \mathfrak{q} \text{ is of type } B_n(i) \text{ with } i \neq n, C_n(i), \text{ or } D_n(i) \\ \varpi_1 & \text{if } \mathfrak{q} \text{ is of type } B_n(n), \end{cases}$$

where ϖ_1 and $\tilde{\varpi}_1$ are the fundamental weights of $\alpha_1 = \varepsilon_1 - \varepsilon_2$ and $\alpha_{i+1} = \varepsilon_{i+1} - \varepsilon_{i+2}$, respectively. As $\mathfrak{l}_{n\gamma}$ acts trivially on both \mathfrak{l}_γ and $\mathfrak{z}(\mathfrak{n})$, the highest weight ν for a constituent $V(\nu) \subset \mathfrak{l}_\gamma \otimes \mathfrak{z}(\mathfrak{n})$ is of the form

$$(6.14) \quad \nu = \sum_{j=1}^{i-1} n_j \varpi_j \quad \text{for } n_j \in \mathbb{Z}_{\geq 0}.$$

If there exists $\beta \in \Delta(\mathfrak{g}(1))$ so that $\nu = \mu + \beta$ then (6.13) and (6.14) imply that $\beta = \nu - \mu$ is of the form

$$(6.15) \quad \begin{cases} (n_1 - 1)\varpi_1 + \sum_{j=2}^{i-1} n_j \varpi_j - \tilde{\varpi}_1 & \text{if } \mathfrak{q} \text{ is of type } B_n(i) \text{ } i \neq n, C_n(i), \text{ or } D_n(i) \\ (n_1 - 1)\varpi_1 + \sum_{j=2}^{i-1} n_j \varpi_j & \text{if } \mathfrak{q} \text{ is of type } B_n(n) \end{cases}$$

for $n_j \in \mathbb{Z}_{\geq 0}$. On the other hand, we observed that the root β must be one of the forms in (6.11). Then observation shows that if β satisfies both (6.11) and (6.15) then β must be

$$\begin{cases} \varpi_1 - \tilde{\varpi}_1 \text{ or } (-\varpi_1 + \varpi_2) - \tilde{\varpi}_1 & \text{if } \mathfrak{q} \text{ is of type } B_n(i) \text{ } i \neq n, C_n(i), \text{ or } D_n(i) \\ \varpi_1 \text{ or } (-\varpi_1 + \varpi_2) & \text{if } \mathfrak{q} \text{ is of type } B_n(n). \end{cases}$$

Therefore $\nu = \mu + \beta$ is $\nu = 2\varpi_1$ or ϖ_2 , which shows that $\nu = 2\varepsilon_1$ or $\varepsilon_1 + \varepsilon_2$. As $\xi_\gamma = \varepsilon_1 - \varepsilon_i$ for \mathfrak{q} of type $B_n(i)$, $C_n(i)$, or $D_n(i)$, Theorem 6.3 shows that both $V(2\varepsilon_1)$ and $V(\varepsilon_1 + \varepsilon_2)$ occur in $\mathfrak{l}_\gamma \otimes \mathfrak{z}(\mathfrak{n})$. Now the assertions follow from the fact that the highest root γ of \mathfrak{g} is $\gamma = \varepsilon_1 + \varepsilon_2$ if \mathfrak{g} is of type B_n or D_n , and $\gamma = 2\varepsilon_1$ if \mathfrak{g} is of type C_n . \square

If \mathfrak{g} is an exceptional algebra then the parabolic subalgebras \mathfrak{q} under consideration are

$$(6.16) \quad E_6(3), E_6(5), E_7(2), E_7(6), E_8(1), \text{ and } F_4(4).$$

Lemma 6.17. *If \mathfrak{q} is of exceptional type as in (6.16) then $V(\xi_\gamma + \gamma_0)$ in Theorem 6.3 is a special constituent.*

Proof. This is done by a direct computation. The roots ϵ_γ in $\Delta(\mathfrak{g}(1))$ so that $\xi_\gamma + \gamma_0 = \mu + \epsilon_\gamma$ are given in Table 5 below. \square

Proposition 6.18. *There exists a unique special constituent in $\mathfrak{l}_\gamma \otimes \mathfrak{z}(\mathfrak{n})$.*

Proof. If \mathfrak{q} is of classical type then this proposition follows from Proposition 6.12. For \mathfrak{q} of exceptional type, by Theorem 6.3, the tensor product $\mathfrak{l}_\gamma \otimes \mathfrak{z}(\mathfrak{n})$ decomposes into

$$\mathfrak{l}_\gamma \otimes \mathfrak{z}(\mathfrak{n}) = V(\xi_\gamma + \gamma) \oplus V(\gamma) \oplus V(\xi_\gamma + \gamma_0)$$

with $\gamma_0 \in \Delta(\mathfrak{n})$ as in Theorem 6.3. Then Lemma 6.10 and Lemma 6.17 show that $V(\xi_\gamma + \gamma_0)$ is the unique special constituent. \square

Since the weight $\epsilon \in \Delta(\mathfrak{g}(1))$ so that $\mu + \epsilon$ is the highest weight of a special constituent will play a role later, we introduce the notation related to ϵ .

Definition 6.19. We denote by ϵ_γ the root contributing to $\mathfrak{g}(1)$ so that $V(\mu + \epsilon_\gamma)$ is the special constituent of $\mathfrak{l}_\gamma \otimes \mathfrak{z}(\mathfrak{n})$. Similarly, we denote by $\epsilon_{n\gamma}$ the root for $\mathfrak{g}(1)$ so that $V(\mu + \epsilon_{n\gamma}) = \mathfrak{l}_{n\gamma} \otimes \mathfrak{z}(\mathfrak{n})$.

We summarize data on the special constituents in Table 2, Table 3, Table 4, and Table 5 below. A dash indicates that no special constituent of the type exists for the case.

By Proposition 6.8, only special constituents could contribute to the construction of the Ω_2 systems. Next we want to show that $\tilde{\tau}_2|_{V^*} \neq 0$ when V is a special constituent. An observation on the highest weights for the special constituents will simplify the argument. We classify them by their highest weights and call them type 1a, type 1b, type 2, and type 3.

Definition 6.20. Let μ be the highest weight for $\mathfrak{g}(1)$, and let $\epsilon = \epsilon_\gamma$ or $\epsilon = \epsilon_{n\gamma}$. (See Definition 6.19.) We say that a special constituent $V(\mu + \epsilon)$ of $\mathfrak{l} \otimes \mathfrak{z}(\mathfrak{n})$ is of

- (1) **type 1a** if $\mu + \epsilon$ is not a root with $\epsilon \neq \mu$ and both μ and ϵ are long roots,
- (2) **type 1b** if $\mu + \epsilon$ is not a root with $\epsilon \neq \mu$ and either μ or ϵ is a short root,
- (3) **type 2** if $\mu + \epsilon = 2\mu$ is not a root, or
- (4) **type 3** if $\mu + \epsilon$ is a root.

Table 6 summarizes the types of special constituents for each parabolic subalgebra \mathfrak{q} . One may want to observe that almost all the special constituents are of type 1a. We regard any roots as long roots, when \mathfrak{g} is simply laced. A dash indicates that no special constituent of the type exists in the case.

Remark 6.21. It is observed from Table 4 and Table 5 that we have $\mu \pm \epsilon \notin \Delta$, unless $V(\mu + \epsilon)$ is of type 3. In particular, if $V(\mu + \epsilon)$ is of type 1a then $\langle \mu, \epsilon \rangle = 0$.

Remark 6.22. Table 6 shows that when $V(\mu + \epsilon)$ is a special constituent of type 1a, \mathfrak{q} is of type $B_n(i)$ ($3 \leq i \leq n-1$), $D_n(i)$, $E_6(3)$, $E_6(5)$, $E_7(2)$, $E_7(6)$, or $E_8(1)$. The data in Appendix A shows

TABLE 2. Highest Weights for Special Constituents (Classical Cases)

Type	$V(\mu + \epsilon_\gamma)$	$V(\mu + \epsilon_{n\gamma})$
$B_n(i)$, $3 \leq i \leq n-2$	$2\epsilon_1$	$\epsilon_1 + \epsilon_2 + \epsilon_{i+1} + \epsilon_{i+2}$
$B_n(n-1)$	$2\epsilon_1$	$\epsilon_1 + \epsilon_2 + \epsilon_n$
$B_n(n)$	$2\epsilon_1$	—
$C_n(i)$, $2 \leq i \leq n-1$	$\epsilon_1 + \epsilon_2$	$2\epsilon_1 + 2\epsilon_{i+1}$
$D_n(i)$, $3 \leq i \leq n-3$	$2\epsilon_1$	$\epsilon_1 + \epsilon_2 + \epsilon_{i+1} + \epsilon_{i+2}$

TABLE 3. Highest Weights for Special Constituents (Exceptional Cases)

Type	$V(\mu + \epsilon_\gamma)$
$E_6(3)$	$\alpha_1 + 2\alpha_2 + 2\alpha_3 + 4\alpha_4 + 3\alpha_5 + 2\alpha_6$
$E_6(5)$	$2\alpha_1 + 2\alpha_2 + 3\alpha_3 + 4\alpha_4 + 2\alpha_5 + \alpha_6$
$E_7(2)$	$2\alpha_1 + 2\alpha_2 + 4\alpha_3 + 5\alpha_4 + 4\alpha_5 + 3\alpha_6 + 2\alpha_7$
$E_7(6)$	$2\alpha_1 + 3\alpha_2 + 4\alpha_3 + 6\alpha_4 + 4\alpha_5 + 2\alpha_6 + \alpha_7$
$E_8(1)$	$2\alpha_1 + 4\alpha_2 + 5\alpha_3 + 8\alpha_4 + 7\alpha_5 + 6\alpha_6 + 4\alpha_7 + 2\alpha_8$
$F_4(4)$	$2\alpha_1 + 4\alpha_2 + 6\alpha_3 + 2\alpha_4$
Type	$V(\mu + \epsilon_{n\gamma})$
$E_6(3)$	$2\alpha_1 + 2\alpha_2 + 2\alpha_3 + 3\alpha_4 + 2\alpha_5 + \alpha_6$
$E_6(5)$	$\alpha_1 + 2\alpha_2 + 2\alpha_3 + 3\alpha_4 + 2\alpha_5 + 2\alpha_6$
$E_7(2)$	—
$E_7(6)$	$2\alpha_1 + 2\alpha_2 + 3\alpha_3 + 4\alpha_4 + 3\alpha_5 + 2\alpha_6 + 2\alpha_7$
$E_8(1)$	—
$F_4(4)$	—

TABLE 4. The Roots μ , ϵ_γ , and $\epsilon_{n\gamma}$ (Classical Cases)

Type	μ	ϵ_γ	$\epsilon_{n\gamma}$
$B_n(i), 3 \leq i \leq n-2$	$\epsilon_1 + \epsilon_{i+1}$	$\epsilon_1 - \epsilon_{i+1}$	$\epsilon_2 + \epsilon_{i+2}$
$B_n(n-1)$	$\epsilon_1 + \epsilon_n$	$\epsilon_1 - \epsilon_n$	ϵ_2
$B_n(n)$	ϵ_1	ϵ_1	—
$C_n(i), 2 \leq i \leq n-1$	$\epsilon_1 + \epsilon_{i+1}$	$\epsilon_2 - \epsilon_{i+1}$	$\epsilon_1 + \epsilon_{i+1}$
$D_n(i), 3 \leq i \leq n-3$	$\epsilon_1 + \epsilon_{i+1}$	$\epsilon_1 - \epsilon_{i+1}$	$\epsilon_2 + \epsilon_{i+2}$

that when \mathfrak{q} is of type $B_n(i)$ for $3 \leq i \leq n-1$, the simple root $\alpha_{\mathfrak{q}} = \epsilon_i - \epsilon_{i+1}$ that parametrizes \mathfrak{q} is a long root and that the set $\Delta(\mathfrak{z}(\mathfrak{n}))$ contains solely long roots. Since we regard any roots as long roots for \mathfrak{g} simply laced, it follows that when $V(\mu + \epsilon)$ is of type 1a, the simple root $\alpha_{\mathfrak{q}}$ and any root $\gamma_j \in \Delta(\mathfrak{z}(\mathfrak{n}))$ are all long roots.

6.3. Computations for Structure Constants. For the rest of this section, we collect technical results on the special constituents and structure constants, so that certain arguments will go smoothly when we find the special values for the Ω_2 systems. The root vectors X_α and the structure constants $N_{\alpha,\beta}$ are normalized as in Section 5.

Lemma 6.23. *Let $V(\mu + \epsilon)$ be a special constituent $\mathfrak{l} \otimes \mathfrak{z}(\mathfrak{n})$ of type 1a, and $\alpha \in \Delta^+(\mathfrak{l})$. If $\epsilon + \alpha \in \Delta$ then $\mu - \alpha \in \Delta$.*

Proof. We show that $\langle \mu, \alpha \rangle > 0$. Since $\mu + \epsilon$ is the highest weight of an irreducible \mathfrak{l} -module, it is $\Delta(\mathfrak{l})$ -dominant. Thus,

$$(6.24) \quad \langle \mu + \epsilon, \alpha \rangle = \langle \mu, \alpha \rangle + \langle \epsilon, \alpha \rangle \geq 0.$$

Observe that, as $\mu + \epsilon$ is of type 1a, ϵ is a long root of \mathfrak{g} . Since $\alpha + \epsilon$ is assumed to be a root, Lemma 4.15 implies that $\langle \alpha, \epsilon^\vee \rangle = -1$; in particular, $\langle \epsilon, \alpha \rangle < 0$. Now, by (6.24), we have

$$\langle \mu, \alpha \rangle \geq -\langle \epsilon, \alpha \rangle > 0.$$

□

TABLE 5. The Roots μ , ϵ_γ , and $\epsilon_{n\gamma}$ (Exceptional Cases)

Type	μ
$E_6(3)$	$\alpha_1 + \alpha_2 + \alpha_3 + 2\alpha_4 + 2\alpha_5 + \alpha_6$
$E_6(5)$	$\alpha_1 + \alpha_2 + 2\alpha_3 + 2\alpha_4 + \alpha_5 + \alpha_6$
$E_7(2)$	$\alpha_1 + \alpha_2 + 2\alpha_3 + 3\alpha_4 + 3\alpha_5 + 2\alpha_6 + \alpha_7$
$E_7(6)$	$\alpha_1 + 2\alpha_2 + 2\alpha_3 + 3\alpha_4 + 2\alpha_5 + \alpha_6 + \alpha_7$
$E_8(1)$	$\alpha_1 + 3\alpha_2 + 3\alpha_3 + 5\alpha_4 + 4\alpha_5 + 3\alpha_6 + 2\alpha_7 + \alpha_8$
$F_4(4)$	$\alpha_1 + 2\alpha_2 + 3\alpha_3 + \alpha_4$
Type	ϵ_γ
$E_6(3)$	$\alpha_2 + \alpha_3 + 2\alpha_4 + \alpha_5 + \alpha_6$
$E_6(5)$	$\alpha_1 + \alpha_2 + \alpha_3 + 2\alpha_4 + \alpha_5$
$E_7(2)$	$\alpha_1 + \alpha_2 + 2\alpha_3 + 2\alpha_4 + \alpha_5 + \alpha_6 + \alpha_7$
$E_7(6)$	$\alpha_1 + \alpha_2 + 2\alpha_3 + 3\alpha_4 + 2\alpha_5 + \alpha_6$
$E_8(1)$	$\alpha_1 + \alpha_2 + 2\alpha_3 + 3\alpha_4 + 3\alpha_5 + 3\alpha_6 + 2\alpha_7 + \alpha_8$
$F_4(4)$	$\alpha_1 + 2\alpha_2 + 3\alpha_3 + \alpha_4$
Type	$\epsilon_{n\gamma}$
$E_6(3)$	$\alpha_1 + \alpha_2 + \alpha_3 + \alpha_4$
$E_6(5)$	$\alpha_2 + \alpha_4 + \alpha_5 + \alpha_6$
$E_7(2)$	—
$E_7(6)$	$\alpha_1 + \alpha_3 + \alpha_4 + \alpha_5 + \alpha_6 + \alpha_7$
$E_8(1)$	—
$F_4(4)$	—

TABLE 6. Types of Special Constituents

Type	$V(\mu + \epsilon_\gamma)$	$V(\mu + \epsilon_{n\gamma})$
$B_n(i), 3 \leq i \leq n-2$	Type 1a	Type 1a
$B_n(n-1)$	Type 1a	Type 1b
$B_n(n)$	Type 2	—
$C_n(i), 2 \leq i \leq n-1$	Type 3	Type 2
$D_n(i), 3 \leq i \leq n-3$	Type 1a	Type 1a
$E_6(3)$	Type 1a	Type 1a
$E_6(5)$	Type 1a	Type 1a
$E_7(2)$	Type 1a	—
$E_7(6)$	Type 1a	Type 1a
$E_8(1)$	Type 1a	—
$F_4(4)$	Type 2	—

Lemma 6.25. *Let $V(\mu + \epsilon)$ be a special constituent of $\mathfrak{l} \otimes \mathfrak{z}(\mathfrak{n})$ of type 1a. If $\alpha \in \Delta^+(\mathfrak{l})$ with $\alpha + \epsilon \in \Delta$ then, for all $\gamma_j \in \Delta(\mathfrak{z}(\mathfrak{n}))$, we have*

$$\text{ad}(X_\mu)\text{ad}(X_{\alpha+\epsilon})X_{-\gamma_j} = 0.$$

Proof. If $(\alpha + \epsilon) - \gamma_j \notin \Delta$ then there is nothing to prove. So we assume that $(\alpha + \epsilon) - \gamma_j \in \Delta$ and $\mu + (\alpha + \epsilon) - \gamma_j \in \Delta$. Since $\mu + \epsilon$ is assumed to be of type 1a, the root μ is long. Lemma 4.15 then implies that

$$(6.26) \quad \langle (\alpha + \epsilon) - \gamma_j, \mu^\vee \rangle = -1.$$

By Remark 6.21, we have $\langle \epsilon, \mu^\vee \rangle = 0$. Thus (6.26) becomes

$$(6.27) \quad \langle \alpha, \mu^\vee \rangle - \langle \gamma_j, \mu^\vee \rangle = -1.$$

Since μ is the highest weight for $\mathfrak{g}(1)$, $\gamma_j \in \Delta(\mathfrak{z}(\mathfrak{n}))$, and $\alpha \in \Delta^+(\mathfrak{l})$, neither $\mu + \alpha$ nor $\gamma_j + \mu$ is a root. Then, as μ is a long root, (6.27) holds if and only if $\langle \alpha, \mu^\vee \rangle = 0$ and $\langle \gamma_j, \mu^\vee \rangle = 1$. On the other hand, since $\alpha + \epsilon$ is a root by hypothesis and by Lemma 6.23, $\mu - \alpha$ is a root. In particular, by Lemma 4.15, $\langle \alpha, \mu^\vee \rangle = 1$. Now we have $\langle \alpha, \mu^\vee \rangle = 1$ and $\langle \alpha, \mu^\vee \rangle = 0$, which is a contradiction. \square

For any $\text{ad}(\mathfrak{h})$ -invariant subspace $W \subset \mathfrak{g}$ and any weight $\nu \in \mathfrak{h}^*$, we write

$$\Delta_\nu(W) = \{\alpha \in \Delta(W) \mid \nu - \alpha \in \Delta\}.$$

In Section 7, we will construct the $\Omega_2|_{V(\mu+\epsilon)^*}$ systems and find their special values, when $V(\mu + \epsilon)$ is of either type 1a or type 2. When we do so, the roots $\beta \in \Delta_{\mu+\epsilon}(\mathfrak{g}(1))$ and $\gamma_j \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))$ will play a role. Therefore, for the rest of this section, we shall show several technical results about those roots, so that certain argument will become simple.

First of all, we need check that $\Delta_{\mu+\epsilon}(\mathfrak{g}(1))$ and $\Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))$ are not empty. It is clear that $\Delta_{\mu+\epsilon}(\mathfrak{g}(1)) \neq \emptyset$, since $\mu, \epsilon \in \Delta_{\mu+\epsilon}(\mathfrak{g}(1))$. Moreover, Lemma 6.28 below shows that when $V(\mu + \epsilon)$ is of type 2, we have $\Delta_{\mu+\epsilon}(\mathfrak{g}(1)) = \{\mu\}$.

Lemma 6.28. *If $V(\mu + \epsilon)$ is a special constituent of $\mathfrak{l} \otimes \mathfrak{z}(\mathfrak{n})$ of type 2 then $\Delta_{\mu+\epsilon}(\mathfrak{g}(1)) = \{\mu\}$.*

Proof. First we claim that μ has the maximum height among the roots $\beta \in \Delta(\mathfrak{g}(1))$. As $\mathfrak{g}(1)$ is the irreducible L -module with highest weight μ , any root $\beta \in \Delta(\mathfrak{g}(1))$ is of the form $\beta = \mu - \sum_{\alpha \in \Pi(\mathfrak{l})} n_\alpha \alpha$ with $n_\alpha \in \mathbb{Z}_{\geq 0}$. Then if $\text{ht}(\mu)$ and $\text{ht}(\beta)$ denote the heights of μ and β , respectively, then

$$\text{ht}(\mu) = \text{ht}(\beta) + \sum_{\alpha \in \Pi(\mathfrak{l})} n_\alpha \geq \text{ht}(\beta).$$

Now as $V(\mu + \epsilon)$ is of type 2, by definition, we have $\mu + \epsilon = 2\mu$. If $\beta \in \Delta_{2\mu}(\mathfrak{g}(1))$ then $2\mu - \beta \in \Delta(\mathfrak{g}(1))$. In particular, the height $\text{ht}(2\mu - \beta)$ satisfies $\text{ht}(\mu) \geq \text{ht}(2\mu - \beta)$. If $\beta = \mu - \sum_{\alpha \in \Pi(\mathfrak{l})} n_\alpha \alpha$ with $n_\alpha \in \mathbb{Z}_{\geq 0}$ then

$$\text{ht}(\mu) \geq \text{ht}(2\mu - \beta) = 2\text{ht}(\mu) - \text{ht}(\beta) = \text{ht}(\mu) + \sum_{\alpha \in \Pi(\mathfrak{l})} n_\alpha.$$

This forces that $\sum_{\alpha \in \Pi(\mathfrak{l})} n_\alpha = 0$. Therefore $\beta = \mu$. \square

Lemma 6.29. *If $V(\mu + \epsilon)$ is a special constituent of $\mathfrak{l} \otimes \mathfrak{z}(\mathfrak{n})$ then $\Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n})) \neq \emptyset$.*

Proof. Observe that the highest weight $\mu + \epsilon$ of $V(\mu + \epsilon) \subset \mathfrak{l} \otimes \mathfrak{z}(\mathfrak{n})$ must be of the form

$$\mu + \epsilon = \begin{cases} \xi_\gamma + \gamma' & \text{if } V(\mu + \epsilon) \subset \mathfrak{l}_\gamma \otimes \mathfrak{z}(\mathfrak{n}) \\ \xi_{n\gamma} + \gamma'' & \text{if } V(\mu + \epsilon) = \mathfrak{l}_{n\gamma} \otimes \mathfrak{z}(\mathfrak{n}) \end{cases}$$

for some $\gamma', \gamma'' \in \Delta(\mathfrak{z}(\mathfrak{n}))$, where ξ_γ and $\xi_{n\gamma}$ are the highest weights for \mathfrak{l}_γ and $\mathfrak{l}_{n\gamma}$, respectively. Then we have $\gamma', \gamma'' \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))$. \square

The following simple technical lemma will simplify an argument in later proofs.

Lemma 6.30. *Let $\alpha, \beta, \delta \in \Delta$ with $\alpha, \beta \neq \delta$. If $\alpha + \beta \notin \Delta$ and $\alpha + \beta - \delta \in \Delta$ then the following hold:*

- (1) $\alpha - \delta, \beta - \delta \in \Delta$, and
- (2) $N_{\beta, \alpha - \delta} N_{\alpha, -\delta} = N_{\alpha, \beta - \delta} N_{\beta, -\delta}$.

Proof. These simply follow from the structure of the complex simple Lie algebras. \square

Lemma 6.31. *Let W be any $\text{ad}(\mathfrak{h})$ -invariant subspace of \mathfrak{g} with the condition $\Delta_{\mu+\epsilon}(W) \setminus \{\mu, \epsilon\} \neq \emptyset$. If $V(\mu + \epsilon)$ is a special constituent of $\mathfrak{l} \otimes \mathfrak{z}(\mathfrak{n})$ of type 1a, type 1b, or type 2 then, for any $\delta \in \Delta_{\mu+\epsilon}(W) \setminus \{\mu, \epsilon\}$, we have $\delta - \mu, \delta - \epsilon \in \Delta$.*

Proof. If $V(\mu + \epsilon)$ is of type 1a, type 1b, or type 2 then, by definition, $\mu + \epsilon$ is not a root. Then this lemma simply follows from Lemma 6.30 \square

Remark 6.32. *A direct observation shows that if $V(\mu + \epsilon)$ is a special constituent of type 1a then $\Delta_{\mu+\epsilon}(\mathfrak{g}(1)) \setminus \{\mu, \epsilon\} \neq \emptyset$.*

Lemma 6.33. *If $V(\mu + \epsilon)$ is a special constituent of $\mathfrak{l} \otimes \mathfrak{z}(\mathfrak{n})$ of type 1a then, for any $\alpha \in \Delta_{\mu+\epsilon}(\mathfrak{g}(1))$ and any $\gamma_j \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))$, we have $\gamma_j - \alpha \in \Delta$.*

Proof. By Lemma 6.30, we have $\gamma_j - \mu, \gamma_j - \epsilon \in \Delta$. So, let $\alpha \neq \mu, \epsilon$. We show that $\langle \gamma_j, \alpha \rangle > 0$. Observe that since $\alpha \in \Delta(\mathfrak{g}(1))$ and $\gamma_j \in \Delta(\mathfrak{z}(\mathfrak{n}))$, we have $\gamma_j + \alpha \notin \Delta$. Thus $\langle \gamma_j, \alpha \rangle \geq 0$. Since $\alpha \in \Delta_{\mu+\epsilon}(\mathfrak{g}(1)) \setminus \{\mu, \epsilon\}$ and $\gamma_j \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))$, by Lemma 6.31, we have $\mu - \alpha, \epsilon - \gamma_j \in \Delta$. Then we first claim that if $\langle \gamma_j, \alpha \rangle = 0$ then $(\mu - \alpha) + (\epsilon - \gamma_j) \in \Delta$. Since $V(\mu + \epsilon)$ is assumed to be of type 1a, both μ and ϵ are long roots. Thus, by Lemma 4.15, $\langle \gamma_j, \mu^\vee \rangle = \langle \alpha, \epsilon^\vee \rangle = 1$; in particular, $\langle \gamma_j, \mu \rangle, \langle \alpha, \epsilon \rangle > 0$. By Remark 6.21, we have $\langle \mu, \epsilon \rangle = 0$. Then,

$$\langle \mu - \alpha, \epsilon - \gamma_j \rangle = -\langle \mu, \gamma_j \rangle - \langle \alpha, \epsilon \rangle < 0.$$

Therefore, as $\mu - \alpha, \epsilon - \gamma_j \in \Delta$, it follows that $(\mu - \alpha) + (\epsilon - \gamma_j) \in \Delta$. On the other hand, since $\langle \mu, \epsilon \rangle = 0$ and $\langle \gamma_j, \alpha \rangle$ is assumed to be 0, we have

$$\begin{aligned} & \|(\mu - \alpha) + (\epsilon - \gamma_j)\|^2 \\ &= \|\mu\|^2 + \|\alpha\|^2 + \|\epsilon\|^2 + \|\gamma_j\|^2 - 2\langle \alpha, \mu \rangle - 2\langle \alpha, \epsilon \rangle - 2\langle \gamma_j, \mu \rangle - 2\langle \gamma_j, \epsilon \rangle. \end{aligned}$$

For $\nu = \alpha, \gamma_j$ and $\zeta = \mu, \epsilon$, by Lemma 4.15, we have $\langle \nu, \zeta^\vee \rangle = 2\langle \nu, \zeta \rangle / \|\zeta\|^2 = 1$, as μ and ϵ are long roots. Therefore, $2\langle \nu, \zeta \rangle = \|\zeta\|^2$, and so,

$$\|(\mu - \alpha) + (\epsilon - \gamma_j)\|^2 = \|\alpha\|^2 + \|\gamma_j\|^2 - \|\mu\|^2 - \|\epsilon\|^2.$$

Since μ and ϵ are assumed to be long roots, this shows that $\|(\mu - \alpha) + (\epsilon - \gamma_j)\|^2 \leq 0$, which contradicts that $(\mu - \alpha) + (\epsilon - \gamma_j)$ is a root. Hence, $\langle \gamma_j, \alpha \rangle > 0$. \square

Lemma 6.34. *If $V(\mu + \epsilon)$ is a special constituent of $\mathfrak{l} \otimes \mathfrak{z}(\mathfrak{n})$ of type 1a or type 2 then, for any $\gamma_j \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))$,*

$$\Delta_{\mu+\epsilon}(\mathfrak{g}(1)) \subset \Delta_{\gamma_j}(\mathfrak{g}(1)).$$

In particular, $\Delta_{\gamma_j}(\mathfrak{g}(1)) \neq \emptyset$ for any $\gamma_j \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))$.

Proof. It is clear that the type 1a case follows from Lemma 6.33. The type 2 case follows from Lemma 6.28 and Lemma 6.31. \square

If $V(\mu + \epsilon)$ is a special constituent of $\mathfrak{l} \otimes \mathfrak{z}(\mathfrak{n})$ then, for $\beta \in \Delta$, we write

$$\theta(\beta) = (\mu + \epsilon) - \beta.$$

Lemma 6.35. *If $V(\mu + \epsilon)$ is a special constituent of $\mathfrak{l} \otimes \mathfrak{z}(\mathfrak{n})$ of type 1a or type 2 then, for any $\gamma_j \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))$,*

$$\Delta_{\theta(\gamma_j)}(\mathfrak{g}(1)) \neq \emptyset.$$

Proof. This simply follows from Lemma 6.30. \square

Lemma 6.36. *If $V(\mu + \epsilon)$ is a special constituent of type 1a or type 2 then*

$$\sum_{\gamma_j \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))} N_{\mu, \epsilon - \gamma_j} N_{-\mu, \gamma_j - \epsilon} N_{\epsilon, -\gamma_j} N_{-\epsilon, \gamma_j} > 0,$$

where $N_{\alpha, \beta}$ are the structure constants for $\alpha, \beta \in \Delta$ defined in Section 5.

Proof. It follows from the normalization (H7) in Section 5 that

$$N_{\mu, \epsilon - \gamma_j} N_{-\mu, \gamma_j - \epsilon} = -\frac{q_{\mu, \epsilon - \gamma_j}(1 + p_{\mu, \epsilon - \gamma_j})}{2} \|\mu\|^2$$

and

$$N_{\epsilon, -\gamma_j} N_{-\epsilon, \gamma_j} = -\frac{q_{\epsilon, -\gamma_j}(1 + p_{\epsilon, -\gamma_j})}{2} \|\epsilon\|^2.$$

In particular, by (5.1) in Section 5, $N_{\mu, \epsilon - \gamma_j} N_{-\mu, \gamma_j - \epsilon} \leq 0$ and $N_{\epsilon, -\gamma_j} N_{-\epsilon, \gamma_j} \leq 0$. By Lemma 6.29 and Lemma 6.31, $\Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n})) \neq \emptyset$ and $\gamma_j - \epsilon \in \Delta$ for any $\gamma_j \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))$. Therefore, for all $\gamma_j \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))$, we have

$$N_{\mu, \epsilon - \gamma_j} N_{-\mu, \gamma_j - \epsilon} N_{\epsilon, -\gamma_j} N_{-\epsilon, \gamma_j} > 0.$$

\square

Lemma 6.37. *If $V(\mu + \epsilon)$ is a special constituent of type 1a or type 2 then, for any $\alpha \in \Delta_{\mu+\epsilon}(\mathfrak{g}(1))$ and any $\gamma_j \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))$, we have the following:*

- (1) $[X_{-\gamma_j}, X_{\alpha - \mu}] = [X_{\theta(\gamma_j)}, X_{\alpha - \mu}] = 0$.
- (2) $N_{\mu - \gamma_j, \alpha - \mu} N_{-(\mu - \gamma_j), -(\alpha - \mu)} = -\frac{\|\mu - \gamma_j\|^2}{2}$.
- (3) $N_{\alpha, -\gamma_j} N_{-\theta(\gamma_j), \theta(\alpha)} = N_{\theta(\alpha), -\gamma_j} N_{-\theta(\gamma_j), \alpha}$.

Proof. To prove (1), we show that $-\gamma_j + \alpha - \mu$ and $\theta(\gamma_j) + \alpha - \mu$ are neither zero nor roots. First of all, if $-\gamma_j + \alpha - \mu = 0$ then $\gamma_j = \mu - \alpha \in \Delta(\mathfrak{l})$, which contradicts that $\gamma_j \in \Delta(\mathfrak{z}(\mathfrak{n}))$. Next, if $\theta(\gamma_j) + \alpha - \mu = 0$ then since $\theta(\gamma_j) + \alpha - \mu = \epsilon + \alpha - \gamma_j$, we would have $\alpha + \epsilon = \gamma_j \in \Delta$. On the other hand, as $V(\mu + \epsilon)$ is assumed to be of type 1a, ϵ is a long root. As $\alpha \in \Delta_{\mu+\epsilon}(\mathfrak{g}(1)) \setminus \{\mu, \epsilon\}$, by Lemma 6.31, we have $\alpha - \epsilon \in \Delta$. Then, by Lemma 4.15, it follows that $\alpha + \epsilon \notin \Delta$, which is a contradiction. To show $\gamma_j + \alpha - \mu$ is not a root, observe that, by Lemma 4.15, we have

$$\langle -\gamma_j + \alpha - \mu, \mu^\vee \rangle = -1 + 1 - 2 = -2.$$

Thus, if $-\gamma_j + \alpha - \mu \in \Delta$ then $(-\gamma_j + \alpha - \mu) + 2\mu$ would be a root. However, since μ is a long root, it is impossible. The fact that $\theta(\gamma_j) + \alpha - \mu \notin \Delta$ can be shown in a similar manner.

By the normalization (H7) in Section 5, to show (2), it suffices to show that $p_{\mu-\gamma_j, \alpha-\mu} = 0$ and $q_{\mu-\gamma_j, \alpha-\mu} = 1$. Observe that, by Lemma 6.33, $(\alpha - \mu) + (\mu - \gamma_j) = \gamma_j - \alpha$ is a root. As $V(\mu + \epsilon)$ is assumed to be of type 1a, μ is a long root. By Remark 6.22, the root γ_j is also a long root. Therefore $\mu - \gamma_j$ is a long root. Now the proposed equality follows from Lemma 4.15.

To show (3), observe that, by the normalization (H3) in Section 5, we have $\kappa(X_\alpha, X_{-\alpha}) = 1$ for all $\alpha \in \Delta$. Thus, $N_{-\theta(\gamma_j), \theta(\alpha)} = \kappa([X_{-\theta(\gamma_j)}, X_{\theta(\alpha)}], X_{\alpha-\gamma_j})$. Then, we have

$$\begin{aligned} N_{-\theta(\gamma_j), \theta(\alpha)} &= \kappa([X_{-\theta(\gamma_j)}, X_{\theta(\alpha)}], X_{\alpha-\gamma_j}) \\ &= \frac{1}{N_{\alpha, -\gamma_j}} \kappa([X_{-\gamma_j}, [X_{-\theta(\gamma_j)}, X_{\theta(\alpha)}]], X_\alpha). \end{aligned}$$

Since $V(\mu + \epsilon)$ is assumed to be of type 1a or type 2, we have $(-\gamma_j) + (-\theta(\gamma_j)) = -(\mu + \epsilon) \notin \Delta$; in particular, $[X_{-\gamma_j}, X_{-\theta(\gamma_j)}] = 0$. Thus, by the Jacobi identity, $[X_{-\gamma_j}, [X_{-\theta(\gamma_j)}, X_{\theta(\alpha)}]] = [X_{-\theta(\gamma_j)}, [X_{-\gamma_j}, X_{\theta(\alpha)}]]$. Hence,

$$\begin{aligned} N_{-\theta(\gamma_j), \theta(\alpha)} &= \frac{1}{N_{\alpha, -\gamma_j}} \kappa([X_{-\gamma_j}, [X_{-\theta(\gamma_j)}, X_{\theta(\alpha)}]], X_\alpha) \\ (6.38) \quad &= \frac{1}{N_{\alpha, -\gamma_j}} N_{-\theta(\gamma_j), \theta(\alpha) - \gamma_j} N_{-\gamma_j, \theta(\alpha)}. \end{aligned}$$

By the normalization (H6), we have $N_{-\theta(\gamma_j), \theta(\alpha) - \gamma_j} = -N_{-\theta(\gamma_j), \alpha}$. Since $N_{-\gamma_j, \theta(\alpha)} = -N_{\theta(\alpha), -\gamma_j}$, it follows from (6.38) that

$$N_{\alpha, -\gamma_j} N_{-\theta(\gamma_j), \theta(\alpha)} = N_{-\theta(\gamma_j), \theta(\alpha) - \gamma_j} N_{-\gamma_j, \theta(\alpha)} = N_{\theta(\alpha), -\gamma_j} N_{-\theta(\gamma_j), \alpha}.$$

□

Lemma 6.39. *If $V(\mu + \epsilon)$ is a special constituent of type 1a then, for any $\alpha \in \Delta_{\mu+\epsilon}(\mathfrak{g}(1)) \setminus \{\mu, \epsilon\}$ and any $\gamma_j \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))$, we have the following:*

- (1) $N_{\alpha, -\gamma_j} N_{\mu, -\alpha} = N_{\mu, -\gamma_j} N_{\alpha-\mu, \mu-\gamma_j}$, and
- (2) $N_{-\theta(\gamma_j), \theta(\alpha)} N_{-\theta(\alpha), \theta(\mu)} = N_{-\theta(\gamma_j), \theta(\mu)} N_{-(\mu-\gamma_j), -(\alpha-\mu)}$.

Proof. It follows from Lemma 6.31 that $\alpha - \mu \in \Delta$. Therefore, we have $X_\alpha = (1/N_{\alpha-\mu, \mu})[X_{\alpha-\mu}, X_\mu]$. Now the assertion (1) follows from the Jacobi identity and the normalization (H6) with Lemma 6.33 and Lemma 6.37 (1). The assertion (2) can be shown similarly. □

Lemma 6.40. *Let \mathfrak{q} be a two-step nilpotent parabolic subalgebra of non-Heisenberg type, listed in (4.9) or (4.10), and $\alpha_{\mathfrak{q}}$ be the simple root that parametrizes the parabolic subalgebra \mathfrak{q} . If $V(\mu + \epsilon)$ is a special constituent of type 1a then, for any $\alpha \in \Delta_{\mu+\epsilon}(\mathfrak{g}(1)) \setminus \{\mu, \epsilon\}$ and any $\gamma_j \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))$,*

$$(6.41) \quad N_{\alpha, -\gamma_j} N_{\mu, -\alpha} N_{-\theta(\gamma_j), \theta(\alpha)} N_{-\theta(\alpha), \theta(\mu)} = N_{\mu, \epsilon - \gamma_j} N_{\epsilon, -\gamma_j} \frac{\|\alpha_{\mathfrak{q}}\|^2}{2}.$$

Proof. By Lemma 6.39, we have

$$\begin{aligned} N_{\alpha, -\gamma_j} N_{\mu, -\alpha} N_{-\theta(\gamma_j), \theta(\alpha)} N_{-\theta(\alpha), \theta(\mu)} &= N_{\mu, -\gamma_j} N_{\alpha-\mu, \mu-\gamma_j} N_{-\theta(\gamma_j), \theta(\mu)} N_{-(\mu-\gamma_j), -(\alpha-\mu)} \\ &= N_{\mu, -\gamma_j} N_{-\theta(\gamma_j), \theta(\mu)} N_{\alpha-\mu, \mu-\gamma_j} N_{-(\mu-\gamma_j), -(\alpha-\mu)} \\ &= N_{\mu, -\gamma_j} N_{-\theta(\gamma_j), \theta(\mu)} \frac{\|\mu - \gamma_j\|^2}{2}. \end{aligned}$$

Note that Lemma 6.37 (2) is applied from line two to line three. Since $-\theta(\gamma_j) + \theta(\mu) + (\mu - \gamma_j) = 0$ with $\theta(\mu) = (\mu + \epsilon) - \mu = \epsilon$, by the normalization (H6), we have $N_{-\theta(\gamma_j), \theta(\mu)} = N_{\epsilon, \mu - \gamma_j}$. By Lemma 6.30 with $\alpha = \mu$, $\beta = \epsilon$, and $\delta = \gamma_j$, it follows that $N_{\epsilon, \mu - \gamma_j} N_{\mu, -\gamma_j} = N_{\mu, \epsilon - \gamma_j} N_{\epsilon, -\gamma_j}$. Therefore,

$$N_{\mu, -\gamma_j} N_{-\theta(\gamma_j), \theta(\mu)} = N_{\mu, -\gamma_j} N_{\epsilon, \mu - \gamma_j} = N_{\mu, \epsilon - \gamma_j} N_{\epsilon, -\gamma_j}.$$

Remark 6.22 shows that γ_j and α_q are long roots, when $V(\mu + \epsilon)$ is of type 1a. Since μ is assumed to be a long root, the root $\mu - \gamma_j$ is a long root. Thus $\|\mu - \gamma_j\|^2 = \|\alpha_q\|^2$. Hence,

$$\begin{aligned} N_{\alpha, -\gamma_j} N_{\mu, -\alpha} N_{-\theta(\gamma_j), \theta(\alpha)} N_{-\theta(\alpha), \theta(\mu)} &= N_{\mu, -\gamma_j} N_{-\theta(\gamma_j), \theta(\mu)} \frac{\|\mu - \gamma_j\|^2}{2} \\ &= N_{\mu, \epsilon - \gamma_j} N_{\epsilon, -\gamma_j} \frac{\|\alpha_q\|^2}{2}. \end{aligned}$$

□

7. THE Ω_2 SYSTEMS

We continue with $\mathfrak{q} = \mathfrak{l} \oplus \mathfrak{g}(1) \oplus \mathfrak{z}(\mathfrak{n})$ a maximal two-step nilpotent parabolic subalgebra of non-Heisenberg type, listed in (4.9) or (4.10). In this section, by using the preliminary results from Section 6, we shall determine the complex parameter $s_2 \in \mathbb{C}$ for the line bundle \mathcal{L}_s so that the Ω_2 systems are conformally invariant on \mathcal{L}_{s_2} . This is done in Theorem 7.16.

7.1. Covariant Map τ_2 . As we have observed in Subsection 3.1, to construct the $\Omega_2|_{V^*}$ system, we use the covariant map τ_2 and the associated L -intertwining operator $\tilde{\tau}_2|_{V^*}$, where V^* is an irreducible constituents of $\mathfrak{l}^* \otimes \mathfrak{z}(\mathfrak{n})^* = \mathfrak{g}(0)^* \otimes \mathfrak{g}(2)^*$. We first show that the covariant map τ_2 is not identically zero, and also that the L -intertwining operators $\tilde{\tau}_2|_{V^*}$ are not identically zero for certain irreducible constituents V . We keep using the normalizations from Section 5.

We start by showing that τ_2 is not identically zero. The covariant map τ_2 is given by

$$\begin{aligned} \tau_2 : \mathfrak{g}(1) &\rightarrow \mathfrak{l} \otimes \mathfrak{z}(\mathfrak{n}) \\ X &\mapsto \frac{1}{2} \operatorname{ad}(X)^2 \omega \end{aligned}$$

with $\omega = \sum_{\gamma_j \in \Delta(\mathfrak{z}(\mathfrak{n}))} X_{-\gamma_j} \otimes X_{\gamma_j}$. The following technical lemma will make a certain argument simpler in later proofs.

Lemma 7.1. *If $V(\mu + \epsilon)$ is a special constituent of type 1a or type 2 then*

$$(7.2) \quad \tau_2(X_\mu + X_\epsilon) = a_{\mu, \epsilon} \operatorname{ad}(X_\mu) \operatorname{ad}(X_\epsilon) \omega,$$

where $a_{\mu, \epsilon} = 1 + \delta_{\mu, \epsilon}$ with $\delta_{\mu, \epsilon}$ the Kronecker delta.

Proof. It is clear that (7.2) holds if $\mu + \epsilon$ is of type 2. Indeed, if $\epsilon = \mu$ then we have

$$\tau_2(2X_\mu) = 4\tau_2(X_\mu) = 2\operatorname{ad}(X_\mu)^2 \omega.$$

If $\mu + \epsilon$ is of type 1a then, by definition, $\mu + \epsilon \notin \Delta$ and both μ and ϵ are long roots. Thus, in the case, $\operatorname{ad}(X_\mu) \operatorname{ad}(X_\epsilon) = \operatorname{ad}(X_\epsilon) \operatorname{ad}(X_\mu)$. Moreover, by Lemma 4.15, we have $\operatorname{ad}(X_\mu)^2 X_{-\gamma_j} = \operatorname{ad}(X_\epsilon)^2 X_{-\gamma_j} = 0$ for any $\gamma_j \in \Delta(\mathfrak{z}(\mathfrak{n}))$. Hence,

$$\tau_2(X_\mu + X_\epsilon) = (1/2)(2 \operatorname{ad}(X_\mu) \operatorname{ad}(X_\epsilon)) \omega = \operatorname{ad}(X_\mu) \operatorname{ad}(X_\epsilon) \omega.$$

□

Proposition 7.3. *Let \mathfrak{q} be a maximal two-step nilpotent parabolic subalgebra of non-Heisenberg type listed in (4.9) or (4.10). Then the covariant map τ_2 is not identically zero.*

Proof. To prove that τ_2 is not identically zero, it suffices to show that there exists a vector $X \in \mathfrak{g}(1)$ so that $\tau_2(X) \neq 0$. Observe that, for each \mathfrak{q} under consideration, $\mathfrak{l} \otimes \mathfrak{z}(\mathfrak{n})$ has at least one special constituent $V(\mu + \epsilon)$ of type 1a or type 2 (see Table 6 in Subsection 6.2). Therefore, $\Delta(\mathfrak{g}(1))$ always contains a root ϵ so that $V(\mu + \epsilon)$ is such a special constituent. Then, to prove this proposition, we show that $\tau_2(X_\mu + X_\epsilon) \neq 0$, where X_μ and X_ϵ are root vectors for μ and ϵ , respectively, with $\mu + \epsilon$ the highest weight for a special constituent of type 1a or type 2.

Let $\mu + \epsilon$ be the highest weight of a special constituent of type 1a or type 2. By Lemma 7.1 we have

$$\tau_2(X_\mu + X_\epsilon) = a_{\mu, \epsilon} \operatorname{ad}(X_\mu) \operatorname{ad}(X_\epsilon) \omega = a_{\mu, \epsilon} \sum_{\gamma_j \in \Delta(\mathfrak{z}(\mathfrak{n}))} \operatorname{ad}(X_\mu) \operatorname{ad}(X_\epsilon) X_{-\gamma_j} \otimes X_{\gamma_j}$$

with $a_{\mu, \epsilon} = 1 + \delta_{\mu, \epsilon}$. If there were a root $\gamma_j \in \Delta(\mathfrak{z}(\mathfrak{n}))$ such that $\epsilon - \gamma_j = -\mu$ then $\mu + \epsilon = \gamma_j \in \Delta$, which contradicts the assumption that $\mu + \epsilon$ is of type 1a or type 2. By Lemma 6.31, if $\mu + \epsilon - \gamma_j \in \Delta$ then $\epsilon - \gamma_j \in \Delta$. Then, for all $\gamma_j \in \Delta(\mathfrak{z}(\mathfrak{n}))$,

$$\operatorname{ad}(X_\mu) \operatorname{ad}(X_\epsilon) X_{-\gamma_j} = \begin{cases} N_{\mu, \epsilon - \gamma_j} N_{\epsilon, -\gamma_j} X_{\mu + \epsilon - \gamma_j} & \text{if } \mu + \epsilon - \gamma_j \in \Delta \\ 0 & \text{otherwise.} \end{cases}$$

Therefore, we have

$$\begin{aligned} \tau_2(X_\mu + X_\epsilon) &= a_{\mu, \epsilon} \sum_{\gamma_j \in \Delta(\mathfrak{z}(\mathfrak{n}))} \operatorname{ad}(X_\mu) \operatorname{ad}(X_\epsilon) X_{-\gamma_j} \otimes X_{\gamma_j} \\ &= a_{\mu, \epsilon} \sum_{\gamma_j \in \Delta_{\mu + \epsilon}(\mathfrak{z}(\mathfrak{n}))} N_{\mu, \epsilon - \gamma_j} N_{\epsilon, -\gamma_j} X_{\mu + \epsilon - \gamma_j} \otimes X_{\gamma_j}. \end{aligned}$$

Since $\{X_{\mu + \epsilon - \gamma_j} \otimes X_{\gamma_j} \mid \gamma_j \in \Delta_{\mu + \epsilon}(\mathfrak{z}(\mathfrak{n}))\}$ is a linearly independent set, this shows that $\tau_2(X_\mu + X_\epsilon) \neq 0$. □

Next we identify irreducible constituent $V(\nu)^*$ so that $\tilde{\tau}_2|_{V(\nu)^*}$ is not identically zero. In Subsection 3.1, we observed that, given an irreducible constituent $V(\nu)^*$, the L -intertwining operator $\tilde{\tau}_2|_{V(\nu)^*} \in \operatorname{Hom}_L(V(\nu)^*, \mathcal{P}^2(\mathfrak{g}(1)))$ is given by

$$(7.4) \quad \tilde{\tau}_2|_{V(\nu)^*}(Y^*)(X) = Y^*(\tau_2(X)),$$

where $\mathcal{P}^2(\mathfrak{g}(1))$ is the space of polynomials on $\mathfrak{g}(1)$ of degree 2. By Proposition 6.8, we know that if $\tilde{\tau}_2|_{V(\nu)^*}$ is not identically zero then $V(\nu)$ is a special constituent of $\mathfrak{l} \otimes \mathfrak{z}(\mathfrak{n})$. We now show that the converse of Proposition 6.8 also holds for special constituents $V(\nu)$ of type 1a or type 2. If $l \in L$ and $Z \in \mathfrak{l}$ then we denote the action of the group and its Lie algebra on $X_\alpha \otimes X_{\gamma_j}$ by $l \cdot (X_\alpha \otimes X_{\gamma_j})$ and $Z \cdot (X_\alpha \otimes X_{\gamma_j})$, respectively.

Proposition 7.5. *If $V(\mu + \epsilon)$ is a special constituent of $\mathfrak{l} \otimes \mathfrak{z}(\mathfrak{n})$ of type 1a or type 2 then the following hold:*

- (1) *The vector $\tau_2(X_\mu + X_\epsilon)$ is a highest weight vector for $V(\mu + \epsilon)$.*

(2) *The L -intertwining operator $\tilde{\tau}_2|_{V(\mu+\epsilon)^*}$ is not identically zero.*

Proof. We have shown that in the proof for Proposition 7.3 that $\tau_2(X_\mu + X_\epsilon) \neq 0$. Moreover, Lemma 7.1 gives that $\tau_2(X_\mu + X_\epsilon) = a_{\mu,\epsilon} \operatorname{ad}(X_\mu) \operatorname{ad}(X_\epsilon) \omega$ with $a_{\mu,\epsilon} = 1 + \delta_{\mu,\epsilon}$. For $l \in L$, we have $l \cdot \omega = \omega$ (see Corollary 3.2) and so

$$l \cdot \tau_2(X_\mu + X_\epsilon) = a_{\mu,\epsilon} \operatorname{ad}(\operatorname{Ad}(l)X_\mu) \operatorname{ad}(\operatorname{Ad}(l)X_\epsilon) \omega.$$

By replacing l by $\exp(tZ)$ with $Z \in \mathfrak{l}$, differentiating, and setting $t = 0$, we obtain

$$(7.6) \quad Z \cdot \tau_2(X_\mu + X_\epsilon) = a_{\mu,\epsilon} (\operatorname{ad}([Z, X_\mu]) \operatorname{ad}(X_\epsilon) + \operatorname{ad}(X_\mu) \operatorname{ad}([Z, X_\epsilon])) \omega.$$

In particular, if $Z = H \in \mathfrak{h}$ in (7.6) then

$$H \cdot \tau_2(X_\mu + X_\epsilon) = (\mu + \epsilon)(H) \tau_2(X_\mu + X_\epsilon).$$

Therefore $\tau_2(X_\mu + X_\epsilon)$ is a weight vector with weight $\mu + \epsilon$. To show that $\tau_2(X_\mu + X_\epsilon)$ is a highest weight vector, we replace Z in (7.6) by X_α with $\alpha \in \Delta^+(\mathfrak{l})$. Since μ is the highest weight for $\mathfrak{g}(1)$, we have

$$X_\alpha \cdot \tau_2(X_\mu + X_\epsilon) = a_{\mu,\epsilon} \operatorname{ad}(X_\mu) \operatorname{ad}([X_\alpha, X_\epsilon]) \omega.$$

If $\mu + \epsilon$ is of type 2 then, as $\epsilon = \mu$ in the case, clearly $X_\alpha \cdot \tau_2(X_\mu + X_\epsilon) = 0$. The case that $\mu + \epsilon$ is of type 1a follows from Lemma 6.25.

To prove the second statement, it is enough to show that there exist $Y^* \in V(\mu + \epsilon)^*$ and $X \in \mathfrak{g}(1)$ so that $\tilde{\tau}_2(Y^*)(X) \neq 0$. Let Y_l^* be a lowest weight vector for $V(\mu + \epsilon)^*$. Observe that if Y_h is a highest weight vector for $V(\mu + \epsilon)$ then $Y_l^*(Y_h) \neq 0$. Since $\tau_2(X_\mu + X_\epsilon)$ is a highest weight vector for $V(\mu + \epsilon)$, we have

$$\tilde{\tau}_2|_{V(\mu+\epsilon)^*}(Y_l^*)(X_\mu + X_\epsilon) = Y_l^*(\tau_2(X_\mu + X_\epsilon)) \neq 0.$$

□

7.2. The $\Omega_2|_{V(\mu+\epsilon)^*}$ Systems. Proposition 7.5 shows that the L -intertwining operator $\tilde{\tau}_2|_{V(\mu+\epsilon)^*}$ is not identically zero, when $V(\mu + \epsilon)$ is a special constituent of $\mathfrak{l} \otimes \mathfrak{z}(\mathfrak{n})$ of type 1a or type 2. We thus construct the $\Omega_2|_{V(\mu+\epsilon)^*}$ system corresponding to irreducible constituents $V(\mu + \epsilon)$ of type 1a or type 2. Here it may be helpful to recall some notation introduced in Subsection 6.3. For any $\operatorname{ad}(\mathfrak{h})$ -invariant subspace $W \subset \mathfrak{g}$ and any weight $\nu \in \mathfrak{h}^*$, we write

$$\Delta_\nu(W) = \{\alpha \in \Delta(W) \mid \nu - \alpha \in \Delta\}.$$

When $V(\mu + \epsilon)$ is a special constituent of $\mathfrak{l} \otimes \mathfrak{z}(\mathfrak{n})$, we write

$$\theta(\beta) = (\mu + \epsilon) - \beta.$$

As indicated in Subsection 3.1, the L -intertwining operator $\tilde{\tau}_2|_{V(\mu+\epsilon)^*}$ yields a system of differential operators. We have denoted such operators by $\Omega_2(Y^*) = \Omega_2|_{V(\mu+\epsilon)^*}(Y^*)$ with $Y^* \in V(\mu + \epsilon)^*$, where $\Omega_2|_{V(\mu+\epsilon)^*} : V(\mu + \epsilon)^* \rightarrow \mathbb{D}(\mathcal{L}_s)^{\bar{\mathfrak{n}}}$ is $\mathcal{U}(\mathfrak{l})$ -equivariant. Because of such equivariance, the system is totally determined, once $\Omega_2(Y_l^*)$ is constructed, where Y_l^* is a lowest weight vector in $V(\mu + \epsilon)^*$.

The first step is to explicitly describe $Y_l^* \in V(\mu + \epsilon)^*$. Observe that we have a non-zero map

$$\begin{aligned}\bar{\tau}_2 : \mathfrak{g}(-1) &\rightarrow \mathfrak{l} \otimes \mathfrak{z}(\mathfrak{n}) \\ \bar{X} &\mapsto \frac{1}{2} \text{ad}(X)^2 \bar{\omega}\end{aligned}$$

with $\bar{\omega} = \sum_{\gamma_t \in \Delta(\mathfrak{z}(\mathfrak{n}))} X_{\gamma_t} \otimes X_{-\gamma_t}$. One checks, as in the proofs for Lemma 3.3 and Proposition 7.3, that $\bar{\tau}_2$ is a non-zero L -equivariant map. Moreover, if $V(\mu + \epsilon)$ is a special constituent of type 1a or type 2 then

$$\bar{\tau}_2(X_{-\mu} + X_{-\epsilon}) = a_{\mu, \epsilon} \text{ad}(X_{-\mu}) \text{ad}(X_{-\epsilon}) \bar{\omega}$$

with $a_{\mu, \epsilon} = 1 + \delta_{\mu, \epsilon}$. Arguing as in Proposition 7.5, we can show that $\bar{\tau}_2(X_{-\mu} + X_{-\epsilon})$ is a lowest weight vector for $V(\mu + \epsilon)^*$ with lowest weight $-\mu - \epsilon$. Thus,

$$(7.7) \quad Y_l^* = \text{ad}(X_{-\mu}) \text{ad}(X_{-\epsilon}) \bar{\omega} = \sum_{\gamma_t \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))} N_{-\mu, \gamma_t - \epsilon} N_{-\epsilon, \gamma_t} X_{-\theta(\gamma_t)} \otimes X_{-\gamma_t}$$

is a lowest weight vector for $V(\mu + \epsilon)^*$. Observe that, by Lemma 6.31, we have $\gamma_t - \epsilon \in \Delta$ for $\gamma_t \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))$. Then, by (7.7), we have

$$\begin{aligned}(7.8) \quad Y_l^*(\tau_2(X)) &= \frac{1}{2} \sum_{\substack{\gamma_t \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n})) \\ \gamma_j \in \Delta(\mathfrak{z}(\mathfrak{n}))}} N_{-\mu, \gamma_t - \epsilon} N_{-\epsilon, \gamma_t} \kappa(X_{-\theta(\gamma_t)}, \text{ad}(X)^2 X_{-\gamma_j}) \kappa(X_{-\gamma_t}, X_{\gamma_j}) \\ &= \frac{1}{2} \sum_{\gamma_t \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))} N_{-\mu, \gamma_t - \epsilon} N_{-\epsilon, \gamma_t} \kappa(X_{-\theta(\gamma_t)}, \text{ad}(X)^2 X_{-\gamma_t}).\end{aligned}$$

Write $X = \sum_{\alpha \in \Delta(\mathfrak{g}(1))} \eta_\alpha X_\alpha$ and let $\gamma_t \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))$. Then,

$$\begin{aligned}\kappa(X_{-\theta(\gamma_t)}, \text{ad}(X)^2 X_{-\gamma_t}) &= \sum_{\alpha, \beta \in \Delta(\mathfrak{g}(1))} \eta_\alpha \eta_\beta \kappa(X_{-\theta(\gamma_t)}, [X_\beta, [X_\alpha, X_{-\gamma_t}]] \\ &= \sum_{\alpha, \beta \in \Delta(\mathfrak{g}(1))} \eta_\alpha \eta_\beta \kappa([X_{-\theta(\gamma_t)}, X_\beta], [X_\alpha, X_{-\gamma_t}]) \\ &= \sum_{\substack{\alpha \in \Delta_{\gamma_t}(\mathfrak{g}(1)) \\ \beta \in \Delta_{\theta(\gamma_t)}(\mathfrak{g}(1))}} \eta_\alpha \eta_\beta N_{\alpha, -\gamma_t} N_{-\theta(\gamma_t), \beta} \kappa(X_{\beta - \theta(\gamma_t)}, X_{\alpha - \gamma_t}).\end{aligned}$$

Observe that, by Lemma 6.34 and Lemma 6.35, the sets $\Delta_{\gamma_t}(\mathfrak{g}(1))$ and $\Delta_{\theta(\gamma_t)}(\mathfrak{g}(1))$ are non-empty. By the normalization (H3) in Section 5, if $\kappa(X_{\beta - \theta(\gamma_t)}, X_{\alpha - \gamma_t}) \neq 0$ then $\beta - \theta(\gamma_t) = \gamma_t - \alpha$. Thus $\kappa(X_{\beta - \theta(\gamma_t)}, X_{\alpha - \gamma_t}) = 0$ unless $\beta = (\mu + \epsilon) - \alpha = \theta(\alpha)$. Therefore,

$$\begin{aligned}(7.9) \quad \kappa(X_{-\theta(\gamma_t)}, \text{ad}(X)^2 X_{-\gamma_t}) &= \sum_{\substack{\alpha \in \Delta_{\gamma_t}(\mathfrak{g}(1)) \\ \beta \in \Delta_{\theta(\gamma_t)}(\mathfrak{g}(1))}} \eta_\alpha \eta_\beta N_{\alpha, -\gamma_t} N_{-\theta(\gamma_t), \beta} \kappa(X_{\beta - \theta(\gamma_t)}, X_{\alpha - \gamma_t}) \\ &= \sum_{\alpha \in \Delta_{\gamma_t}(\mathfrak{g}(1)) \cap \Delta_{\mu+\epsilon}(\mathfrak{g}(1))} N_{\alpha, -\gamma_t} N_{-\theta(\gamma_t), \theta(\alpha)} \eta_\alpha \eta_{\theta(\alpha)} \\ &= \sum_{\alpha \in \Delta_{\mu+\epsilon}(\mathfrak{g}(1))} N_{\alpha, -\gamma_t} N_{-\theta(\gamma_t), \theta(\alpha)} \kappa(X, X_{-\alpha}) \kappa(X, X_{-\theta(\alpha)}).\end{aligned}$$

Lemma 6.34 is used in line three to show that $\Delta_{\gamma_t}(\mathfrak{g}(1)) \cap \Delta_{\mu+\epsilon}(\mathfrak{g}(1)) = \Delta_{\mu+\epsilon}(\mathfrak{g}(1))$. Hence, by (7.8) and (7.9), $\tilde{\tau}_2|_{V(\mu+\epsilon)^*}(Y_l^*)(X) = Y_l^*(\tau_2(X))$ is

$$\begin{aligned} & \tilde{\tau}_2|_{V(\mu+\epsilon)^*}(Y_l^*)(X) \\ &= \frac{1}{2} \sum_{\substack{\alpha \in \Delta_{\mu+\epsilon}(\mathfrak{g}(1)) \\ \gamma_t \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))}} (N_{-\mu, \gamma_t-\epsilon} N_{-\epsilon, \gamma_t})(N_{\alpha, -\gamma_t} N_{-\theta(\gamma_t), \theta(\alpha)}) \kappa(X, X_{-\alpha}) \kappa(X, X_{-\theta(\alpha)}). \end{aligned}$$

Now, via the composition of maps

$$V(\mu + \epsilon)^* \xrightarrow{\tilde{\tau}_2|_{V(\mu+\epsilon)^*}} \mathcal{P}^2(\mathfrak{g}(1)) \rightarrow \text{Sym}^2(\mathfrak{g}(-1)) \xrightarrow{\sigma} \mathcal{U}(\bar{\mathfrak{n}}) \xrightarrow{R} \mathbb{D}(\mathcal{L}_s)^{\bar{\mathfrak{n}}},$$

for $Y_l^* \in V(\mu + \epsilon)^*$, the second-order differential operator $\Omega_2(Y_l^*) \in \mathbb{D}(\mathcal{L}_s)^{\bar{\mathfrak{n}}}$ is given by

$$\Omega_2(Y_l^*) = \frac{1}{2} \sum_{\substack{\alpha \in \Delta_{\mu+\epsilon}(\mathfrak{g}(1)) \\ \gamma_t \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))}} (N_{-\mu, \gamma_t-\epsilon} N_{-\epsilon, \gamma_t})(N_{\alpha, -\gamma_t} N_{-\theta(\gamma_t), \theta(\alpha)})^{\top} R(X_{-\alpha}) R(X_{-\theta(\alpha)})^{\top},$$

where $^{\top}ab^{\top} = (1/2)(ab + ba)$. By Lemma 6.37 (3), no symmetrization is needed. Therefore we obtain

$$(7.10) \quad \Omega_2(Y_l^*) = \frac{1}{2} \sum_{\substack{\alpha \in \Delta_{\mu+\epsilon}(\mathfrak{g}(1)) \\ \gamma_t \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))}} (N_{-\mu, \gamma_t-\epsilon} N_{-\epsilon, \gamma_t})(N_{\alpha, -\gamma_t} N_{-\theta(\gamma_t), \theta(\alpha)}) R(X_{-\alpha}) R(X_{-\theta(\alpha)}).$$

7.3. Special Values of the $\Omega_2|_{V(\mu+\epsilon)^*}$ Systems. Now we determine the special values of the line bundle \mathcal{L}_s for which the $\Omega_2|_{V(\mu+\epsilon)^*}$ system is conformally invariant, under the assumption that $V(\mu + \epsilon)$ is a special constituent of type 1a or type 2.

Choose a basis $\{Y_1^*, \dots, Y_n^*\}$ for $V(\mu + \epsilon)^*$. To show that the list of differential operators $\Omega_2(Y_1^*), \dots, \Omega_2(Y_n^*)$ is conformally invariant on the bundle \mathcal{L}_s , we need to prove that in $\mathbb{D}(\mathcal{L}_s)^{\bar{\mathfrak{n}}}$,

$$(7.11) \quad [\pi_s(X), \Omega_2(Y_i^*)] \in \text{span}_{C^\infty(\bar{N}_0)} \{\Omega_2(Y_1^*), \dots, \Omega_2(Y_n^*)\}$$

for all $X \in \mathfrak{g}$ and all i . By Proposition 2.16, (7.11) holds if

$$(7.12) \quad [\pi_s(X), \Omega_2(Y_i^*)]_e \in \text{span}_{\mathbb{C}} \{\Omega_2(Y_1^*)_e, \dots, \Omega_2(Y_n^*)_e\}$$

holds for all $X \in \mathfrak{g}$ and all i . Here, for $D \in \mathbb{D}(\mathcal{L}_s)$, $D_{\bar{n}}$ denotes the linear functional $f \mapsto (D \bullet f)(\bar{n})$ for $f \in C^\infty(\bar{N}_0, \mathbb{C}_{\chi^s})$. We show that a simplification of (7.12) implies (7.11).

Proposition 7.13. *Let $V(\mu + \epsilon)^*$ be the dual module of a special constituent $V(\mu + \epsilon)$ of $\mathfrak{l} \otimes \mathfrak{z}(\mathfrak{n})$ with respect to the Killing form. Suppose that the operator $\Omega_2|_{V(\mu+\epsilon)^*} : V(\mu + \epsilon)^* \rightarrow \mathbb{D}(\mathcal{L}_s)^{\bar{\mathfrak{n}}}$ is non-zero. If X_h is a highest weight vector for $\mathfrak{g}(1)$ and if we have*

$$[\pi_s(X_h), \Omega_2(Y_l^*)]_e \in \text{span}_{\mathbb{C}} \{\Omega_2(Y_1^*)_e, \dots, \Omega_2(Y_n^*)_e\}$$

for a lowest weight vector Y_l^ and a basis $\{Y_1^*, \dots, Y_n^*\}$ for $V(\mu + \epsilon)^*$ then the $\Omega_2|_{V(\mu+\epsilon)^*}$ system is a conformally invariant system.*

Proof. By Remark 3.8, the $\Omega_k|_{V(\mu+\epsilon)^*}$ system satisfies the condition (S1) of Definition 2.3. We need to prove that (7.12) holds for all $X \in \mathfrak{g} = \bar{\mathfrak{n}} \oplus \mathfrak{l} \oplus \mathfrak{n}$. Note that, by definition, we have $\Omega_2(Y_i^*) \in \mathbb{D}(\mathcal{L}_s)^{\bar{\mathfrak{n}}}$. Hence (7.12) holds for $X \in \bar{\mathfrak{n}}$ trivially. The L_0 -equivariance of $\Omega_2|_{V(\mu+\epsilon)^*}$ shows

that (7.12) holds for $X \in \mathfrak{l}$. Furthermore, Lemma 3.12 established (7.12) when $X \in \mathfrak{g}(1)$. Now we handle the case when $X \in \mathfrak{z}(\mathfrak{n})$.

If $X \in \mathfrak{z}(\mathfrak{n})$ then, since $\mathfrak{z}(\mathfrak{n}) = [\mathfrak{g}(1), \mathfrak{g}(1)]$, it is of the form $X = [X_1, X_2]$ for some $X_1, X_2 \in \mathfrak{g}(1)$. Then, by the Jacobi identity, we have

$$[\pi_s(X), \Omega_2(Y_i^*)] = [\pi_s(X_1), [\pi_s(X_2), \Omega_2(Y_i^*)]] - [\pi_s(X_2), [\pi_s(X_1), \Omega_2(Y_i^*)]].$$

By (2.9), we have $\pi_s(X_j)_e = 0$ for $j = 1, 2$. It follows from Lemma 3.12 that for $j = 1, 2$ and all i , we have

$$[\pi_s(X_j), \Omega_2(Y^*)]_e \in \text{span}_{\mathbb{C}}\{\Omega_2(Y_1^*)_e, \dots, \Omega_2(Y_n^*)_e\}.$$

Therefore, by Lemma 3.9,

$$\begin{aligned} [\pi_s(X), \Omega_2(Y_i^*)]_e &= [\pi_s(X_1), [\pi_s(X_2), \Omega_2(Y_i^*)]]_e - [\pi_s(X_2), [\pi_s(X_1), \Omega_2(Y_i^*)]]_e \\ &\in \text{span}_{\mathbb{C}}\{\Omega_2(Y_1^*)_e, \dots, \Omega_2(Y_k^*)_e\}. \end{aligned}$$

□

Proposition 7.14. *If μ is the highest weight for $\mathfrak{g}(1)$ and $\alpha, \beta \in \Delta(\mathfrak{g}(1))$ then*

$$\begin{aligned} &[\pi_s(X_\mu), R(X_{-\alpha})R(X_{-\beta})]_e \\ &= R([X_\mu, X_{-\alpha}], X_{-\beta})_e - s\lambda_{\mathfrak{q}}([X_\mu, X_{-\alpha}])R(X_{-\beta})_e - s\lambda_{\mathfrak{q}}([X_\mu, X_{-\beta}])R(X_{-\alpha})_e. \end{aligned}$$

Proof. This simply follows by substituting $Y = X_\mu$, $X_1 = X_{-\alpha}$, and $X_2 = X_{-\beta}$ in Proposition 2.17, and evaluating at $\bar{n} = e$. □

If $V(\mu + \epsilon)$ is a special constituent of $\mathfrak{l} \otimes \mathfrak{z}(\mathfrak{n})$ of type 1a or type 2 then we write

$$(7.15) \quad C(\mu, \epsilon) := \sum_{\gamma_t \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))} N_{\mu, \epsilon - \gamma_t} N_{-\mu, \gamma_t - \epsilon} N_{\epsilon, -\gamma_t} N_{-\epsilon, \gamma_t}.$$

By Lemma 6.36, we have $C(\mu, \epsilon) \neq 0$.

Theorem 7.16. *Let \mathfrak{g} be a complex simple Lie algebra and \mathfrak{q} be a maximal two-step nilpotent parabolic subalgebra of non-Heisenberg type, listed in (4.9) or (4.10). If Y_l^* is the lowest weight vector defined in (7.7) for the dual module $V(\mu + \epsilon)^*$ of a special constituent $V(\mu + \epsilon)$ of type 1a or type 2, and if $\alpha_{\mathfrak{q}}$ is the simple root that determines \mathfrak{q} then the following hold:*

(1) *If $V(\mu + \epsilon)$ is of type 1a then*

$$(7.17) \quad [\pi_s(X_\mu), \Omega_2(Y_l^*)]_e = -\frac{\|\alpha_{\mathfrak{q}}\|^2}{2} C(\mu, \epsilon) (s - s_2) R(X_{-\epsilon})_e,$$

with $s_2 = \frac{|\Delta_{\mu+\epsilon}(\mathfrak{g}(1))|}{2} - 1$, where $|\Delta_{\mu+\epsilon}(\mathfrak{g}(1))|$ is the number of elements in $\Delta_{\mu+\epsilon}(\mathfrak{g}(1))$.

(2) *If $V(\mu + \epsilon)$ is of type 2 then*

$$(7.18) \quad [\pi_s(X_\mu), \Omega_2(Y_l^*)]_e = -\frac{\|\alpha_{\mathfrak{q}}\|^2}{2} C(\mu, \mu) (s + 1) R(X_{-\mu})_e.$$

Proof. We start by showing that (7.17) holds. It follows from (7.10) that

$$(7.19) \quad \begin{aligned} & [\pi_s(X_\mu), \Omega_2(Y_l^*)]_e \\ &= \frac{1}{2} \sum_{\substack{\alpha \in \Delta_{\mu+\epsilon}(\mathfrak{g}(1)) \\ \gamma_t \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))}} (N_{-\mu, \gamma_t - \epsilon} N_{-\epsilon, \gamma_t}) (N_{\alpha, -\gamma_t} N_{-\theta(\gamma_t), \theta(\alpha)}) [\pi_s(X_\mu), R(X_{-\alpha}) R(X_{-\theta(\alpha)})]_e. \end{aligned}$$

We use Proposition 7.14 to compute $[\pi_s(X_\mu), R(X_{-\alpha}) R(X_{-\theta(\alpha)})]_e$. This is

$$\begin{aligned} & [\pi_s(X_\mu), R(X_{-\alpha}) R(X_{-\theta(\alpha)})]_e = \\ & R([X_\mu, X_{-\alpha}], X_{-\theta(\alpha)})_e - s\lambda_{\mathfrak{q}}([X_\mu, X_{-\alpha}]) R(X_{-\theta(\alpha)})_e - s\lambda_{\mathfrak{q}}([X_\mu, X_{-\theta(\alpha)}]) R(X_{-\alpha})_e. \end{aligned}$$

We consider the contributions from each term in (7.19), separately. Recall here that, as we defined in Section 4.2, our parabolic subalgebra \mathfrak{q} is parametrized by the simple root $\alpha_{\mathfrak{q}} \in \Pi$ and that $\lambda_{\mathfrak{q}}$ is the fundamental weight for $\alpha_{\mathfrak{q}}$.

First we study the contribution from the second term. It is

$$T_2 = -\frac{s}{2} \sum_{\substack{\alpha \in \Delta_{\mu+\epsilon}(\mathfrak{g}(1)) \\ \gamma_t \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))}} (N_{-\mu, \gamma_t - \epsilon} N_{-\epsilon, \gamma_t}) (N_{\alpha, -\gamma_t} N_{-\theta(\gamma_t), \theta(\alpha)}) \lambda_{\mathfrak{q}}([X_\mu, X_{-\alpha}]) R(X_{-\theta(\alpha)})_e.$$

As $\mathfrak{g}(1)$ is the 1-eigenspace of $\text{ad}(H_{\mathfrak{q}})$ with $H_{\mathfrak{q}}$ defined in (4.4), the set $\Delta(\mathfrak{g}(1))$ is $\Delta(\mathfrak{g}(1)) = \{\beta \in \Delta \mid 2\langle \lambda_{\mathfrak{q}}, \beta \rangle / \|\alpha_{\mathfrak{q}}\|^2 = 1\}$. Therefore, by the normalization (H4) in Section 5, for $\beta \in \Delta(\mathfrak{g}(1))$, we have $\lambda_{\mathfrak{q}}(H_\beta) = \langle \lambda_{\mathfrak{q}}, \beta \rangle = \|\alpha_{\mathfrak{q}}\|^2 / 2$. Thus,

$$(7.20) \quad \lambda_{\mathfrak{q}}([X_\mu, X_{-\alpha}]) = \frac{\|\alpha_{\mathfrak{q}}\|^2}{2} \delta_{\alpha, \mu}$$

with $\delta_{\alpha, \mu}$ the Kronecker delta. So the contribution from this term is

$$\begin{aligned} T_2 &= -\frac{s}{2} \sum_{\substack{\alpha \in \Delta_{\mu+\epsilon}(\mathfrak{g}(1)) \\ \gamma_t \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))}} (N_{-\mu, \gamma_t - \epsilon} N_{-\epsilon, \gamma_t}) (N_{\alpha, -\gamma_t} N_{-\theta(\gamma_t), \theta(\alpha)}) \lambda_{\mathfrak{q}}([X_\mu, X_{-\alpha}]) R(X_{-\theta(\alpha)})_e \\ &= -\frac{s\|\alpha_{\mathfrak{q}}\|^2}{4} \sum_{\gamma_t \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))} (N_{-\mu, \gamma_t - \epsilon} N_{-\epsilon, \gamma_t}) (N_{\mu, -\gamma_t} N_{-\theta(\gamma_t), \theta(\mu)}) R(X_{-\theta(\mu)})_e \\ &= -\frac{s\|\alpha_{\mathfrak{q}}\|^2}{4} \sum_{\gamma_t \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))} (N_{-\mu, \gamma_t - \epsilon} N_{-\epsilon, \gamma_t}) (N_{\mu, -\gamma_t} N_{\epsilon, \mu - \gamma_t}) R(X_{-\epsilon})_e. \end{aligned}$$

We showed in Lemma 6.30 that $N_{\mu, -\gamma_t} N_{\epsilon, \mu - \gamma_t} = N_{\mu, \epsilon - \gamma_t} N_{\epsilon, -\gamma_t}$. Hence,

$$\begin{aligned} T_2 &= -\frac{s\|\alpha_{\mathfrak{q}}\|^2}{4} \sum_{\gamma_t \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))} (N_{-\mu, \gamma_t - \epsilon} N_{-\epsilon, \gamma_t}) (N_{\mu, -\gamma_t} N_{\epsilon, \mu - \gamma_t}) R(X_{-\epsilon})_e \\ &= -\frac{s\|\alpha_{\mathfrak{q}}\|^2}{4} \sum_{\gamma_t \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))} (N_{-\mu, \gamma_t - \epsilon} N_{-\epsilon, \gamma_t}) (N_{\mu, \epsilon - \gamma_t} N_{\epsilon, -\gamma_t}) R(X_{-\epsilon})_e \\ &= -\frac{s\|\alpha_{\mathfrak{q}}\|^2}{4} \sum_{\gamma_t \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))} (N_{\mu, \epsilon - \gamma_t} N_{-\mu, \gamma_t - \epsilon} N_{\epsilon, -\gamma_t} N_{-\epsilon, \gamma_t}) R(X_{-\epsilon})_e \\ &= -\frac{s\|\alpha_{\mathfrak{q}}\|^2}{4} C(\mu, \epsilon) R(X_{-\epsilon})_e. \end{aligned}$$

The same argument shows that the contribution from the third term is

$$\begin{aligned} T_3 &= -\frac{s}{2} \sum_{\substack{\alpha \in \Delta_{\mu+\epsilon}(\mathfrak{g}(1)) \\ \gamma_t \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))}} (N_{-\mu, \gamma_t - \epsilon} N_{-\epsilon, \gamma_t}) (N_{\alpha, -\gamma_t} N_{-\theta(\gamma_t), \theta(\alpha)}) \lambda_{\mathfrak{q}}([X_{\mu}, X_{-\theta(\alpha)}]) R(X_{-\alpha})_e \\ &= -\frac{s \|\alpha_{\mathfrak{q}}\|^2}{4} C(\mu, \epsilon) R(X_{-\epsilon})_e. \end{aligned}$$

Now we consider the contribution from the first term. It is

$$T_1 = \frac{1}{2} \sum_{\substack{\alpha \in \Delta_{\mu+\epsilon}(\mathfrak{g}(1)) \\ \gamma_t \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))}} (N_{-\mu, \gamma_t - \epsilon} N_{-\epsilon, \gamma_t}) (N_{\alpha, -\gamma_t} N_{-\theta(\gamma_t), \theta(\alpha)}) R([X_{\mu}, X_{-\alpha}], X_{-\theta(\alpha)})_e.$$

We claim that if $\alpha = \epsilon$ or μ then $[[X_{\mu}, X_{-\alpha}], X_{-\theta(\alpha)}] = 0$, where $\theta(\alpha)$ denotes $\theta(\alpha) = (\mu + \epsilon) - \alpha$. If $\alpha = \epsilon$ then, by Remark 6.21, $[X_{\mu}, X_{-\alpha}] = [X_{\mu}, X_{-\epsilon}] = 0$. If $\alpha = \mu$ then

$$[[X_{\mu}, X_{-\mu}], X_{-\theta(\mu)}] = [[X_{\mu}, X_{-\mu}], X_{-\epsilon}] = \epsilon(H_{\mu})X_{-\epsilon} = 0.$$

Note that Remark 6.21 is applied to obtain $\epsilon(H_{\mu}) = \langle \epsilon, \mu \rangle = 0$. Moreover, by Remark 6.32, we have $\Delta_{\mu+\epsilon}(\mathfrak{g}(1)) \setminus \{\mu, \epsilon\} \neq \emptyset$. The contribution from T_1 is

$$\begin{aligned} T_1 &= \frac{1}{2} \sum_{\substack{\alpha \in \Delta_{\mu+\epsilon}(\mathfrak{g}(1)) \\ \gamma_t \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))}} (N_{-\mu, \gamma_t - \epsilon} N_{-\epsilon, \gamma_t}) (N_{\alpha, -\gamma_t} N_{-\theta(\gamma_t), \theta(\alpha)}) R([X_{\mu}, X_{-\alpha}], X_{-\theta(\alpha)})_e \\ &= \frac{1}{2} \sum_{\substack{\alpha \in \Delta_{\mu+\epsilon}(\mathfrak{g}(1)) \setminus \{\mu, \epsilon\} \\ \gamma_t \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))}} (N_{-\mu, \gamma_t - \epsilon} N_{-\epsilon, \gamma_t}) (N_{\alpha, -\gamma_t} N_{-\theta(\gamma_t), \theta(\alpha)}) R([X_{\mu}, X_{-\alpha}], X_{-\theta(\alpha)})_e \\ &= \frac{1}{2} \sum_{\substack{\alpha \in \Delta_{\mu+\epsilon}(\mathfrak{g}(1)) \setminus \{\mu, \epsilon\} \\ \gamma_t \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))}} (N_{-\mu, \gamma_t - \epsilon} N_{-\epsilon, \gamma_t}) (N_{\alpha, -\gamma_t} N_{-\theta(\gamma_t), \theta(\alpha)}) (N_{\mu, -\alpha} N_{-\mu - \alpha, -\theta(\alpha)}) R(X_{-\epsilon})_e \\ &= \frac{1}{2} \sum_{\substack{\alpha \in \Delta_{\mu+\epsilon}(\mathfrak{g}(1)) \setminus \{\mu, \epsilon\} \\ \gamma_t \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))}} (N_{-\mu, \gamma_t - \epsilon} N_{-\epsilon, \gamma_t}) (N_{\alpha, -\gamma_t} N_{-\theta(\gamma_t), \theta(\alpha)}) (N_{\mu, -\alpha} N_{-\theta(\alpha), \theta(\mu)}) R(X_{-\epsilon})_e \\ &= \frac{1}{2} \sum_{\substack{\alpha \in \Delta_{\mu+\epsilon}(\mathfrak{g}(1)) \setminus \{\mu, \epsilon\} \\ \gamma_t \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))}} (N_{-\mu, \gamma_t - \epsilon} N_{-\epsilon, \gamma_t}) (N_{\alpha, -\gamma_t} N_{\mu, -\alpha} N_{-\theta(\gamma_t), \theta(\alpha)} N_{-\theta(\alpha), \theta(\mu)}) R(X_{-\epsilon})_e. \end{aligned}$$

Note that, from line three to line four, we use that $N_{\mu - \alpha, -\theta(\alpha)} = N_{-\theta(\alpha), \theta(\mu)}$, as $(\mu - \alpha) + (-\theta(\alpha)) + \theta(\mu) = 0$. (See the normalization (H6) in Section 5.) By Lemma 6.40, we have

$$N_{\alpha, -\gamma_t} N_{\mu, -\alpha} N_{-\theta(\gamma_t), \theta(\alpha)} N_{-\theta(\alpha), \theta(\mu)} = N_{\mu, \epsilon - \gamma_t} N_{\epsilon, -\gamma_t} \frac{\|\alpha_{\mathfrak{q}}\|^2}{2}.$$

Therefore,

$$\begin{aligned}
T_1 &= \frac{1}{2} \sum_{\substack{\alpha \in \Delta_{\mu+\epsilon}(\mathfrak{g}(1)) \setminus \{\mu, \epsilon\} \\ \gamma_t \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))}} (N_{-\mu, \gamma_t - \epsilon} N_{-\epsilon, \gamma_t}) (N_{\alpha, -\gamma_t} N_{\mu, -\alpha} N_{-\theta(\gamma_t), \theta(\alpha)} N_{-\theta(\alpha), \theta(\mu)}) R(X_{-\epsilon})_e \\
&= \frac{\|\alpha_q\|^2}{4} \sum_{\substack{\alpha \in \Delta_{\mu+\epsilon}(\mathfrak{g}(1)) \setminus \{\mu, \epsilon\} \\ \gamma_t \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))}} (N_{-\mu, \gamma_t - \epsilon} N_{-\epsilon, \gamma_t}) (N_{\mu, \epsilon - \gamma_t} N_{\epsilon, -\gamma_t}) R(X_{-\epsilon})_e \\
&= \frac{\|\alpha_q\|^2}{4} \sum_{\substack{\alpha \in \Delta_{\mu+\epsilon}(\mathfrak{g}(1)) \setminus \{\mu, \epsilon\} \\ \gamma_t \in \Delta_{\mu+\epsilon}(\mathfrak{z}(\mathfrak{n}))}} (N_{\mu, \epsilon - \gamma_t} N_{-\mu, \gamma_t - \epsilon} N_{\epsilon, -\gamma_t} N_{-\epsilon, \gamma_t}) R(X_{-\epsilon})_e \\
&= \frac{\|\alpha_q\|^2}{4} C(\mu, \epsilon) \sum_{\alpha \in \Delta_{\mu+\epsilon}(\mathfrak{g}(1)) \setminus \{\mu, \epsilon\}} R(X_{-\epsilon})_e \\
&= \frac{\|\alpha_q\|^2}{4} C(\mu, \epsilon) (|\Delta_{\mu+\epsilon}(\mathfrak{g}(1))| - 2) R(X_{-\epsilon})_e.
\end{aligned}$$

Hence, we obtain

$$\begin{aligned}
[\pi_s(X_\mu), \Omega_2(Y_l^*)]_e &= T_1 + T_2 + T_3 \\
&= -\frac{\|\alpha_q\|^2}{2} C(\mu, \epsilon) \left(s - \left(\frac{|\Delta_{\mu+\epsilon}(\mathfrak{g}(1))|}{2} - 1 \right) \right) R(X_{-\epsilon})_e.
\end{aligned}$$

Now we are going to prove the equation (7.18). If $V(\mu + \epsilon)$ is of type 2 then $\mu + \epsilon = 2\mu$; in particular, $\theta(\mu) = (2\mu) - \mu = \mu$. By Lemma 6.28, $\Delta_{2\mu}(\mathfrak{g}(1)) = \{\mu\}$. Thus, (7.10) becomes

$$\begin{aligned}
\Omega_2(Y_l^*) &= \frac{1}{2} \sum_{\substack{\alpha \in \Delta_{2\mu}(\mathfrak{g}(1)) \\ \gamma_t \in \Delta_{2\mu}(\mathfrak{z}(\mathfrak{n}))}} (N_{-\mu, \gamma_t - \epsilon} N_{-\epsilon, \gamma_t}) (N_{\alpha, -\gamma_t} N_{-\theta(\gamma_t), \theta(\alpha)}) R(X_{-\alpha}) R(X_{-\theta(\alpha)}) \\
&= \frac{1}{2} \sum_{\gamma_t \in \Delta_{2\mu}(\mathfrak{z}(\mathfrak{n}))} (N_{-\mu, \gamma_t - \mu} N_{-\mu, \gamma_t}) (N_{\mu, -\gamma_t} N_{-\theta(\gamma_t), \theta(\mu)}) R(X_{-\mu}) R(X_{-\theta(\mu)}) \\
(7.21) \quad &= \frac{1}{2} \sum_{\gamma_t \in \Delta_{2\mu}(\mathfrak{z}(\mathfrak{n}))} (N_{-\mu, \gamma_t - \mu} N_{-\mu, \gamma_t}) (N_{\mu, -\gamma_t} N_{-\theta(\gamma_t), \mu}) R(X_{-\mu})^2.
\end{aligned}$$

Since $(-\theta(\gamma_t)) + \mu + (\mu - \gamma_t) = 0$, we have $N_{-\theta(\gamma_t), \mu} = N_{\mu, \mu - \gamma_t}$. Thus,

$$\begin{aligned}
&\frac{1}{2} \sum_{\gamma_t \in \Delta_{2\mu}(\mathfrak{z}(\mathfrak{n}))} (N_{-\mu, \gamma_t - \mu} N_{-\mu, \gamma_t}) (N_{\mu, -\gamma_t} N_{-\theta(\gamma_t), \mu}) R(X_{-\mu})^2 \\
&= \frac{1}{2} \sum_{\gamma_t \in \Delta_{2\mu}(\mathfrak{z}(\mathfrak{n}))} (N_{-\mu, \gamma_t - \mu} N_{-\mu, \gamma_t}) (N_{\mu, -\gamma_t} N_{\mu, \mu - \gamma_t}) R(X_{-\mu})^2 \\
&= \frac{1}{2} \sum_{\gamma_t \in \Delta_{2\mu}(\mathfrak{z}(\mathfrak{n}))} (N_{\mu, \mu - \gamma_t} N_{-\mu, \gamma_t - \mu} N_{\mu, -\gamma_t} N_{-\mu, \gamma_t}) R(X_{-\mu})^2 \\
(7.22) \quad &= \frac{1}{2} C(\mu, \mu) R(X_{-\mu})^2.
\end{aligned}$$

Therefore,

$$[\pi_s(X_\mu), \Omega_2(Y_l^*)]_e = \frac{1}{2} C(\mu, \mu) [\pi_s(X_\mu), R(X_{-\mu})^2]_e.$$

It follows from (7.20) that $\lambda_{\mathfrak{q}}([X_{\mu}, X_{-\mu}]) = \|\alpha_{\mathfrak{q}}\|^2/2$. Then, by Proposition 7.14 with $\alpha = \beta = \mu$, we have

$$\begin{aligned} [\pi_s(X_{\mu}), R(X_{-\mu})^2]_e &= R([X_{\mu}, X_{-\mu}], X_{-\mu})_e - 2s\lambda_{\mathfrak{q}}([X_{\mu}, X_{-\mu}])R(X_{-\mu})_e \\ &= -\mu(H_{\mu})R(X_{-\mu})_e - 2s \cdot \frac{\|\alpha_{\mathfrak{q}}\|^2}{2}R(X_{-\mu})_e \\ &= -(s\|\alpha_{\mathfrak{q}}\|^2 + \|\mu\|^2)R(X_{-\mu})_e. \end{aligned}$$

Observe that Table 6 in Subsection 6.2 shows that a special constituent of type 2 occurs only when \mathfrak{q} is of type $B_n(n)$, type $C_n(i)$ or $F_4(4)$. Appendix A shows that when \mathfrak{q} is of these types, we have $\|\mu\|^2 = \|\alpha_{\mathfrak{q}}\|^2$. Therefore,

$$[\pi_s(X_{\mu}), R(X_{-\mu})^2]_e = -(s\|\alpha_{\mathfrak{q}}\|^2 + \|\mu\|^2)R(X_{-\mu})_e = -\|\alpha_{\mathfrak{q}}\|^2(s+1)R(X_{-\mu})_e.$$

Hence, we obtain

$$\begin{aligned} [\pi_s(X_{\mu}), \Omega_2(Y_l^*)]_e &= \frac{1}{2}C(\mu, \mu)[\pi_s(X_{\mu}), R(X_{-\mu})^2]_e \\ &= -\frac{\|\alpha_{\mathfrak{q}}\|^2}{2}C(\mu, \mu)(s+1)R(X_{-\mu})_e. \end{aligned}$$

□

To emphasize the fundamental weight $\lambda_{\mathfrak{q}}$, we write $\mathcal{L}(s\lambda_{\mathfrak{q}})$ for the line bundle \mathcal{L}_s . Now, by combining Proposition 7.13 and Theorem 7.16, we conclude the following.

Corollary 7.23. *Under the same hypotheses in Theorem 7.16, we have:*

- (1) *If $V(\mu + \epsilon)^*$ is of type 1a then the $\Omega_2|_{V(\mu + \epsilon)^*}$ system is conformally invariant on the line bundle $\mathcal{L}(s_2\lambda_{\mathfrak{q}})$, where s_2 is the constant given in Theorem 7.16.*
- (2) *If $V(\mu + \epsilon)^*$ is of type 2 then the $\Omega_2|_{V(\mu + \epsilon)^*}$ system is conformally invariant on the line bundle $\mathcal{L}(-\lambda_{\mathfrak{q}})$.*

Proof. This corollary follows from Proposition 7.13 and Theorem 7.16. □

As we defined in Definition 6.19, we denote by $V(\mu + \epsilon_{\gamma})$ the special constituent of $\mathfrak{l} \otimes \mathfrak{z}(\mathfrak{n})$ so that $V(\mu + \epsilon_{\gamma}) \subset \mathfrak{l}_{\gamma} \otimes \mathfrak{z}(\mathfrak{n})$, and denote by $V(\mu + \epsilon_{n\gamma})$ the special constituent so that $V(\mu + \epsilon_{n\gamma}) = \mathfrak{l}_{n\gamma} \otimes \mathfrak{z}(\mathfrak{n})$. See Table 6 in Subsection 6.2 for the types of $V(\mu + \epsilon_{\gamma})$ and $V(\mu + \epsilon_{n\gamma})$ for each case. Table 7 below summarizes the line bundles $\mathcal{L}(s_0\lambda_{\mathfrak{q}})$ on which the Ω_2 systems are conformally invariant. Here, a dash indicates that there does not exist the special constituent $V(\mu + \epsilon_{n\gamma})$. When \mathfrak{q} is of type $B_n(n-1)$, the special constituent $V(\mu + \epsilon_{n\gamma})$ is of type 1b, and when \mathfrak{q} is of type $C_n(i)$, the special constituent $V(\mu + \epsilon_{\gamma})$ is of type 3. Therefore, we put a question mark for these cases in the table.

APPENDIX A. MISCELLANEOUS DATA

This appendix summarizes the miscellenious data for the maximal two-step nilpotent parabolic subalgebras $\mathfrak{q} = \mathfrak{l} \oplus \mathfrak{g}(1) \oplus \mathfrak{z}(\mathfrak{n})$ of non-Heisenberg type shown in (4.9) and (4.10) in Section 4. For each case we give the deleted Dynkin diagram of \mathfrak{q} , the subgraphs for \mathfrak{l}_{γ} and $\mathfrak{l}_{n\gamma}$, the simple root α_{γ} that is not orthogonal to the highest root for \mathfrak{g} , the highest weights for $\mathfrak{g}(1)$ and $\mathfrak{z}(\mathfrak{n})$, and the highest roots for \mathfrak{l}_{γ} and $\mathfrak{l}_{n\gamma}$. For the definition for the deleted Dynkin diagram see Subsection 4.1.

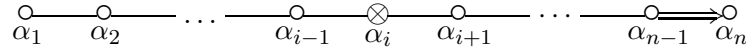
TABLE 7. Line Bundles with Special Values

Parabolic subalgebra \mathfrak{q}	$\Omega_2 _{V(\mu+\epsilon_\gamma)^*}$	$\Omega_2 _{V(\mu+\epsilon_{n\gamma})^*}$
$B_n(i), 3 \leq i \leq n-2$	$\mathcal{L}((n-i-\frac{1}{2})\lambda_i)$	$\mathcal{L}(\lambda_i)$
$B_n(n-1)$	$\mathcal{L}(\frac{1}{2}\lambda_{n-1})$?
$B_n(n)$	$\mathcal{L}(-\lambda_n)$	—
$C_n(i), 2 \leq i \leq n-1$?	$\mathcal{L}(-\lambda_i)$
$D_n(i), 3 \leq i \leq n-3$	$\mathcal{L}((n-i-1)\lambda_i)$	$\mathcal{L}(\lambda_i)$
$E_6(3)$	$\mathcal{L}(\lambda_3)$	$\mathcal{L}(2\lambda_3)$
$E_6(5)$	$\mathcal{L}(\lambda_5)$	$\mathcal{L}(2\lambda_5)$
$E_7(2)$	$\mathcal{L}(2\lambda_2)$	—
$E_7(6)$	$\mathcal{L}(\lambda_6)$	$\mathcal{L}(3\lambda_6)$
$E_8(1)$	$\mathcal{L}(3\lambda_1)$	—
$F_4(4)$	$\mathcal{L}(-\lambda_4)$	—

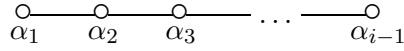
Subsection 4.2 describes about the subspaces $\mathfrak{g}(1)$ and $\mathfrak{z}(\mathfrak{n})$. The definitions for the simple ideals \mathfrak{l}_γ and $\mathfrak{l}_{n\gamma}$ of \mathfrak{l} are given in Subsection 4.3. For classical algebras the sets of roots contributing to for \mathfrak{l} -modules $\mathfrak{g}(1)$, $\mathfrak{z}(\mathfrak{n})$, \mathfrak{l}_γ , and $\mathfrak{l}_{n\gamma}$ are given in the standard realization of the roots.

$$\S B_n(i), 3 \leq i \leq n-2$$

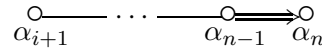
(1) The deleted Dynkin diagram:



(2) The subgraph for \mathfrak{l}_γ :



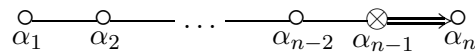
(3) The subgraph for $\mathfrak{l}_{n\gamma}$:



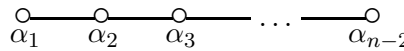
We have $\alpha_\gamma = \alpha_2$. The highest weight μ and the set of roots $\Delta(\mathfrak{g}(1))$ for $\mathfrak{g}(1)$ are $\mu = \varepsilon_1 + \varepsilon_{i+1}$ and $\Delta(\mathfrak{g}(1)) = \{\varepsilon_j \pm \varepsilon_k \mid 1 \leq j \leq i \text{ and } i+1 \leq k \leq n\} \cup \{\varepsilon_j \mid 1 \leq j \leq i\}$. The highest weight γ and the set of roots $\Delta(\mathfrak{z}(\mathfrak{n}))$ for $\mathfrak{z}(\mathfrak{n})$ are $\gamma = \varepsilon_1 + \varepsilon_2$ and $\Delta(\mathfrak{z}(\mathfrak{n})) = \{\varepsilon_j + \varepsilon_k \mid 1 \leq j < k \leq i\}$. The highest root ξ_γ and the set of positive roots $\Delta^+(\mathfrak{l}_\gamma)$ for \mathfrak{l}_γ are $\xi_\gamma = \varepsilon_1 - \varepsilon_i$ and $\Delta^+(\mathfrak{l}_\gamma) = \{\varepsilon_j - \varepsilon_k \mid 1 \leq j < k \leq i\}$. The highest root $\xi_{n\gamma}$ and the set of positive roots $\Delta^+(\mathfrak{l}_{n\gamma})$ for $\mathfrak{l}_{n\gamma}$ are $\xi_{n\gamma} = \varepsilon_{i+1} + \varepsilon_{i+2}$ and $\Delta^+(\mathfrak{l}_{n\gamma}) = \{\varepsilon_j \pm \varepsilon_k \mid i+1 \leq j < k \leq n\} \cup \{\varepsilon_j \mid i+1 \leq j \leq n\}$.

$$\S B_n(n-1)$$

(1) The deleted Dynkin diagram:



(2) The subgraph for \mathfrak{l}_γ :



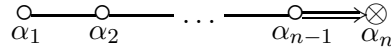
(3) The subgraph for $\mathfrak{l}_{n\gamma}$:



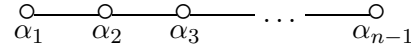
We have $\alpha_\gamma = \alpha_2$. The highest weight μ and the set of weights $\Delta(\mathfrak{g}(1))$ for $\mathfrak{g}(1)$ are $\mu = \varepsilon_1 + \varepsilon_n$ and $\Delta(\mathfrak{g}(1)) = \{\varepsilon_j \pm \varepsilon_n \mid 1 \leq j \leq n-1\} \cup \{\varepsilon_j \mid 1 \leq j \leq n-1\}$. The highest weight γ and the set of weights $\Delta(\mathfrak{z}(\mathfrak{n}))$ for $\mathfrak{z}(\mathfrak{n})$ are $\gamma = \varepsilon_1 + \varepsilon_2$ and $\Delta(\mathfrak{z}(\mathfrak{n})) = \{\varepsilon_j + \varepsilon_k \mid 1 \leq j < k \leq n-1\}$. The highest root ξ_γ and the set of positive roots $\Delta^+(\mathfrak{l}_\gamma)$ for \mathfrak{l}_γ are $\xi_\gamma = \varepsilon_1 - \varepsilon_{n-1}$ and $\Delta^+(\mathfrak{l}_\gamma) = \{\varepsilon_j - \varepsilon_k \mid 1 \leq j < k \leq n-1\}$. The highest root $\xi_{n\gamma}$ and the set of positive roots $\Delta^+(\mathfrak{l}_{n\gamma})$ for $\mathfrak{l}_{n\gamma}$ are $\xi_{n\gamma} = \varepsilon_n$ and $\Delta^+(\mathfrak{l}_{n\gamma}) = \{\varepsilon_n\}$.

§B_n(n)

(1) The deleted Dynkin diagram:



(2) The subgraph for \mathfrak{l}_γ :

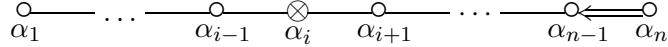


(3) No subgraph for $\mathfrak{l}_{n\gamma}$ ($\mathfrak{l}_{n\gamma} = \{0\}$)

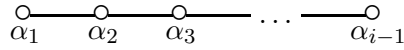
We have $\alpha_\gamma = \alpha_2$. The highest weight μ and the set of weights $\Delta(\mathfrak{g}(1))$ are $\mu = \varepsilon_1$ and $\Delta(\mathfrak{g}(1)) = \{\varepsilon_j \mid 1 \leq j \leq n\}$. The highest weight γ and the set of weights $\Delta(\mathfrak{z}(\mathfrak{n}))$ for $\mathfrak{z}(\mathfrak{n})$ are $\gamma = \varepsilon_1 + \varepsilon_2$ and $\Delta(\mathfrak{z}(\mathfrak{n})) = \{\varepsilon_j + \varepsilon_k \mid 1 \leq j < k \leq n\}$. The highest root ξ_γ and the set of positive roots for \mathfrak{l}_γ are $\xi_\gamma = \varepsilon_1 - \varepsilon_n$ and $\Delta^+(\mathfrak{l}_\gamma) = \{\varepsilon_j - \varepsilon_k \mid 1 \leq j < k \leq n\}$.

§C_n(i), $2 \leq i \leq n-1$

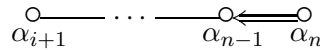
(1) The deleted Dynkin diagram:



(2) The subgraph for \mathfrak{l}_γ :



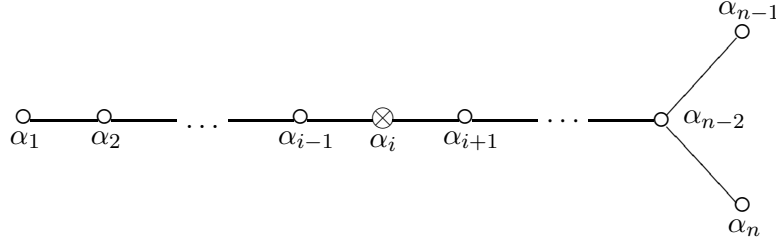
(3) The subgraph for $\mathfrak{l}_{n\gamma}$:



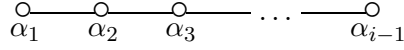
We have $\alpha_\gamma = \alpha_1$. The highest weight μ and the set of weights $\Delta(\mathfrak{g}(1))$ for $\mathfrak{g}(1)$ are $\mu = \varepsilon_1 + \varepsilon_{i+1}$ and $\Delta(\mathfrak{g}(1)) = \{\varepsilon_j \pm \varepsilon_k \mid 1 \leq j \leq i \text{ and } i+1 \leq k \leq n\}$. The highest weight γ and the set of weights $\Delta(\mathfrak{z}(\mathfrak{n}))$ for $\mathfrak{z}(\mathfrak{n})$ are $\gamma = 2\varepsilon_1$ and $\Delta(\mathfrak{z}(\mathfrak{n})) = \{\varepsilon_j + \varepsilon_k \mid 1 \leq j < k \leq i\} \cup \{2\varepsilon_j \mid 1 \leq j \leq i\}$. The highest root ξ_γ and the set of positive roots $\Delta^+(\mathfrak{l}_\gamma)$ for \mathfrak{l}_γ are $\xi_\gamma = \varepsilon_1 - \varepsilon_i$ and $\Delta^+(\mathfrak{l}_\gamma) = \{\varepsilon_j - \varepsilon_k \mid 1 \leq j < k \leq i\}$. The highest root $\xi_{n\gamma}$ and the set of positive roots $\Delta(\mathfrak{l}_{n\gamma})$ for $\mathfrak{l}_{n\gamma}$ are $\xi_{n\gamma} = 2\varepsilon_{i+1}$ and $\Delta^+(\mathfrak{l}_{n\gamma}) = \{\varepsilon_j \pm \varepsilon_k \mid i+1 \leq j < k \leq n\} \cup \{2\varepsilon_j \mid i+1 \leq j \leq n\}$.

§D_n(i), $3 \leq i \leq n-3$

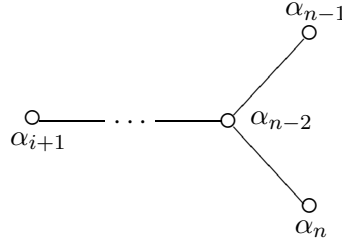
(1) The deleted Dynkin diagram:



(2) The subgraph for \mathfrak{l}_γ :



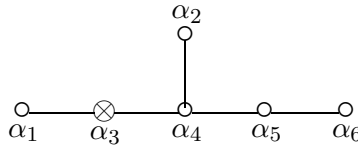
(3) The subgraph for $\mathfrak{l}_{n\gamma}$:



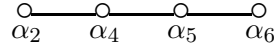
We have $\alpha_\gamma = \alpha_2$. The highest weight μ and the set of weights $\Delta(\mathfrak{g}(1))$ for $\mathfrak{g}(1)$ are $\mu = \varepsilon_1 + \varepsilon_{i+1}$ and $\Delta(\mathfrak{g}(1)) = \{\varepsilon_j \pm \varepsilon_k \mid 1 \leq j \leq i \text{ and } i+1 \leq k \leq n\}$. The highest weight γ and the set of weights $\Delta(\mathfrak{z}(\mathfrak{n}))$ for $\mathfrak{z}(\mathfrak{n})$ are $\gamma = \varepsilon_1 + \varepsilon_2$ and $\Delta(\mathfrak{z}(\mathfrak{n})) = \{\varepsilon_j + \varepsilon_k \mid 1 \leq j < k \leq i\}$. The highest root ξ_γ and the set of positive roots $\Delta^+(\mathfrak{l}_\gamma)$ for \mathfrak{l}_γ are $\xi_\gamma = \varepsilon_1 - \varepsilon_i$ and $\Delta^+(\mathfrak{l}_\gamma) = \{\varepsilon_j - \varepsilon_k \mid 1 \leq j < k \leq i\}$. The highest root $\xi_{n\gamma}$ and the set of positive roots $\Delta^+(\mathfrak{l}_{n\gamma})$ for $\mathfrak{l}_{n\gamma}$ are $\xi_{n\gamma} = \varepsilon_{i+1} + \varepsilon_{i+2}$ and $\Delta^+(\mathfrak{l}_{n\gamma}) = \{\varepsilon_j \pm \varepsilon_k \mid i+1 \leq j < k \leq n\}$.

§E₆(3)

(1) The deleted Dynkin diagram:



(2) The subgraph for \mathfrak{l}_γ :



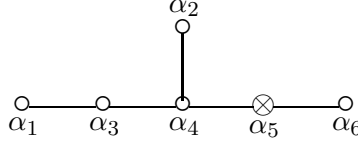
(3) The subgraph for $\mathfrak{l}_{n\gamma}$:



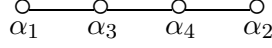
We have $\alpha_\gamma = \alpha_2$. The highest weight μ for $\mathfrak{g}(1)$ is $\mu = \alpha_1 + \alpha_2 + \alpha_3 + 2\alpha_4 + 2\alpha_5 + \alpha_6$. The highest weight γ for $\mathfrak{z}(\mathfrak{n})$ is $\gamma = \alpha_1 + 2\alpha_2 + 2\alpha_3 + 3\alpha_4 + 2\alpha_5 + \alpha_6$. The highest root ξ_γ for \mathfrak{l}_γ is $\xi_\gamma = \alpha_2 + \alpha_4 + \alpha_5 + \alpha_6$. The highest root $\xi_{n\gamma}$ for $\mathfrak{l}_{n\gamma}$ is $\xi_{n\gamma} = \alpha_1$.

§E₆(5)

(1) The deleted Dynkin diagram:



(2) The subgraph for \mathfrak{l}_γ :



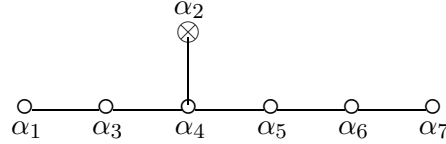
(3) The subgraph for $\mathfrak{l}_{n\gamma}$:



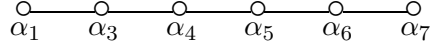
We have $\alpha_\gamma = \alpha_2$. The highest weight μ for $\mathfrak{g}(1)$ is $\mu = \alpha_1 + \alpha_2 + 2\alpha_3 + 2\alpha_4 + \alpha_5 + \alpha_6$. The highest weight γ for $\mathfrak{z}(\mathfrak{n})$ is $\gamma = \alpha_1 + 2\alpha_2 + 2\alpha_3 + 3\alpha_4 + 2\alpha_5 + \alpha_6$. The highest weight ξ_γ for \mathfrak{l}_γ is $\xi_\gamma = \alpha_1 + \alpha_2 + \alpha_3 + \alpha_4$. The highest weight $\xi_{n\gamma}$ for $\mathfrak{l}_{n\gamma}$ is $\xi_{n\gamma} = \alpha_6$.

§E₇(2)

(1) The deleted Dynkin diagram:



(2) The subgraph for \mathfrak{l}_γ :

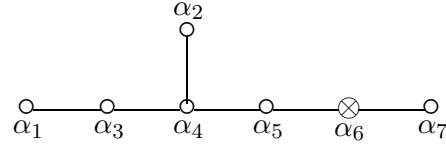


(3) No subgraph for $\mathfrak{l}_{n\gamma}$ ($\mathfrak{l}_{n\gamma} = \{0\}$)

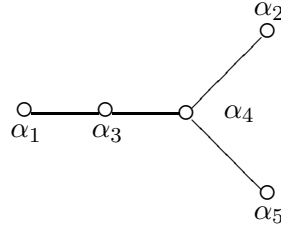
We have $\alpha_\gamma = \alpha_1$. The highest weight μ for $\mathfrak{g}(1)$ is $\mu = \alpha_1 + \alpha_2 + 2\alpha_3 + 3\alpha_4 + 3\alpha_5 + 2\alpha_6 + \alpha_7$. The highest weight γ for $\mathfrak{z}(\mathfrak{n})$ is $\gamma = 2\alpha_1 + 2\alpha_2 + 3\alpha_3 + 4\alpha_4 + 3\alpha_5 + 2\alpha_6 + \alpha_7$. The highest root ξ_γ for \mathfrak{l}_γ is $\xi_\gamma = \alpha_1 + \alpha_3 + \alpha_4 + \alpha_5 + \alpha_6 + \alpha_7$.

§E₇(6)

(1) The deleted Dynkin diagram:



(2) The subgraph for \mathfrak{l}_γ :



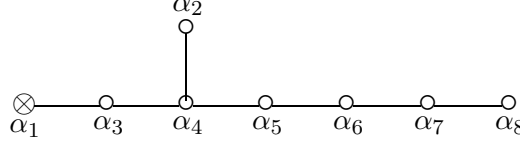
(3) The subgraph for $\mathfrak{l}_{n\gamma}$:



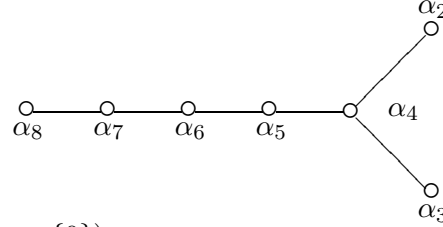
We have $\alpha_\gamma = \alpha_1$. The highest weight μ for $\mathfrak{g}(1)$ is $\mu = \alpha_1 + 2\alpha_2 + 2\alpha_3 + 3\alpha_4 + 2\alpha_5 + \alpha_6 + \alpha_7$. The highest weight γ for $\mathfrak{z}(\mathfrak{n})$ is $\gamma = 2\alpha_1 + 2\alpha_2 + 3\alpha_3 + 4\alpha_4 + 3\alpha_5 + 2\alpha_6 + \alpha_7$. The highest root ξ_γ for \mathfrak{l}_γ is $\xi_\gamma = \alpha_1 + \alpha_2 + 2\alpha_3 + 2\alpha_4 + \alpha_5$. The highest root $\xi_{n\gamma}$ for $\mathfrak{l}_{n\gamma}$ is $\xi_{n\gamma} = \alpha_7$.

§E₈(1)

(1) The deleted Dynkin diagram:



(2) The subgraph for \mathfrak{l}_γ :

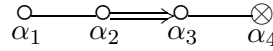


(3) No subgraph for $\mathfrak{l}_{n\gamma}$ ($\mathfrak{l}_{n\gamma} = \{0\}$)

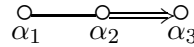
We have $\alpha_\gamma = \alpha_8$. The highest weight μ for $\mathfrak{g}(1)$ is $\mu = \alpha_1 + 3\alpha_2 + 3\alpha_3 + 5\alpha_4 + 4\alpha_5 + 3\alpha_6 + 2\alpha_7 + \alpha_8$. The highest weight γ for $\mathfrak{z}(\mathfrak{n})$ is $\gamma = 2\alpha_1 + 3\alpha_2 + 4\alpha_3 + 6\alpha_4 + 5\alpha_5 + 4\alpha_6 + 3\alpha_7 + 2\alpha_8$. The highest root ξ_γ for \mathfrak{l}_γ is $\xi_\gamma = \alpha_2 + \alpha_3 + 2\alpha_4 + 2\alpha_5 + 2\alpha_6 + 2\alpha_7 + \alpha_8$.

§F₄(4)

(1) The deleted Dynkin diagram:



(2) The subgraph for \mathfrak{l}_γ :



(3) No subgraph for $\mathfrak{l}_{n\gamma}$ ($\mathfrak{l}_{n\gamma} = \{0\}$)

We have $\alpha_\gamma = \alpha_1$. The highest weight μ for $\mathfrak{g}(1)$ is $\mu = \alpha_1 + 2\alpha_2 + 3\alpha_3 + \alpha_4$. The highest weight γ for $\mathfrak{z}(\mathfrak{n})$ is $\gamma = 2\alpha_1 + 3\alpha_2 + 4\alpha_3 + 2\alpha_4$. The highest root for ξ_γ for \mathfrak{l}_γ is $\xi_\gamma = \alpha_1 + 2\alpha_2 + 2\alpha_3$.

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